GEORGE A. HAGEDORN

ALAIN JOYE

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Annales de l'I. H. P., section A, tome 68, nº 1 (1998), p. 85-134 http://www.numdam.org/item?id=AIHPA_1998_68_1_85_0

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Landau-Zener transitions through small electronic eigenvalue gaps in the Born-Oppenheimer approximation

by

George A. HAGEDORN * and Alain JOYE **, ***

Department of Mathematics and Center for Statistical Mechanics and Mathematical Physics, Virginia Polytechnic Institude and State University, Blacksburg, Virginia 24061-0123, U.S.A.

ABSTRACT. – We study the propagation of molecular wave packets through the simplest two types of avoided crossings of electronic energy levels in a limit where the gap between the eigenvalues shrinks as the nuclear masses are increased. For these types of avoided crossings, the electron energy levels essentially depend on only one of the nuclear configuration parameters, as is the case for all diatomic molecules. We find that the transition probabilities are of order 1 and are determined by the Landau-Zener formula. © Elsevier, Paris.

Key words: Born-Oppenheimer approximation, avoided crossings, molecular dynamics, Landau-Zener transitions, adiabatic approximations.

RÉSUMÉ. – Nous étudions la propagation de paquets d'ondes moléculaires au travers des deux types les plus simples de croisements évités de niveaux d'énergie électroniques, dans la limite où le gap entre ces valeurs propres décroît lorsque les masses nucléaires augmentent. Pour ces types de croisements évités, les niveaux d'énergie électroniques ne dépendent essentiellement que d'un seul des paramètres de configuration nucléaire,

Annales de l'Institut Henri Poincaré - Physique théorique - 0246-0211 Vol. 68/98/01/© Elsevier, Paris

^{*} Partially supported by National Science Foundation Grant DMS-9403401

^{**} Supported by Fonds National Suisse de la Recherche Scientifique, Grant 8220-037200.

^{***} Permanent address: Centre de Physique Théorique, CNRS Marseille, Luminy Case 907,

¹³²⁸⁸ Marseille Cedex 9, France and PHYMAT, Université de Toulon et du Var, B.P.132, 83957 La Garde Cedex. France.

comme c'est le cas pour toute molécule diatomique. Nous observons que les probabilités de transition sont d'ordre un et sont déterminées par la formule de Landau-Zener. © Elsevier, Paris.

1. INTRODUCTION

Because it is not practical to solve the time-dependent Schrödinger equation directly, the time-dependent Born-Oppenheimer approximation is a principal tool for studying molecular dynamics. This approximation makes use of the smallness of the parameter ϵ , where ϵ^4 is the ratio of the mass of an electron to the average of the masses of the nuclei.

In the standard time-dependent Born-Oppenheimer approximation, the electrons and nuclei are treated separately, but their motions are coupled. The electrons move much faster than the nuclei, and they quickly adjust their motion in response to the relatively slow nuclear motion. The electrons remain approximately in a quantum mechanical bound state as though the nuclei were at fixed classical positions. This is the *adiabatic approximation* for the electrons. The motion of the nuclei is accurately described by the *semiclassical approximation* because the nuclei have large masses. The electronic and nuclear motions are coupled because the energy level of the electronic bound state depends on the positions of the nuclei, and the electronic energy level plays the role of an effective potential for the semiclassical dynamics of the nuclei.

This physical intuition is the basis for rigorous asymptotic expansions of solutions to the molecular time-dependent Schrödinger equation [6], [8], [11], [15], [16]. However, the validity of the approximation is dependent upon the assumption that the electron energy level of interest is well isolated from the rest of the spectrum of the electronic Hamiltonian.

Readers interested in the mathematical literature concerning the validity of Born-Oppenheimer approximations should consult [2-4], [6], [8-13], [15-17], [20], [32], [34-36], [38-42].

The assumption that the electron energy level of interest is isolated from the rest of the spectrum can break down in various ways. For example, two electron energy levels may cross one another for some nuclear configurations. The effects of such crossings on molecular propagation have been studied recently in generic minimal multiplicity situations [13], [16]. [17]. Another situation where the standard approximation breaks down is an "avoided crossing," where two electron energy levels approach close to one another, but do not actually cross. Generic avoided crossings of energy levels that have the minimal multiplicity allowed by the symmetry group have been classified, and normal forms for the electron Hamiltonian near these avoided crossings have been determined [18]. In [18] it is shown that there are six distinct types of these avoided crossings.

In this paper we study molecular propagation through the simplest two types of avoided crossings described in [18]. In these types of crossings, the electron energy levels essentially depend on only one parameter in the nuclear configuration space. In practice, this occurs for diatomic molecules, where the electron energy levels depend only on the distance between the nuclei because of rotational symmetry. There are two types of such avoided crossings because of the possible presence of time reversing operators in the symmetry group of the electron Hamiltonian. The details of the situation dictate whether minimal multiplicity energy levels are of multiplicity 1 (Type 1 Avoided Crossings) or multiplicity 2 (Type 2 Avoided Crossings).

Our main result is the determination of what happens when a standard time-dependent Born-Oppenheimer molecular wave packet propagates through one of these avoided crossings if the gap size is on the order of ϵ . Using matched asymptotic expansions we explicitly compute approximate solutions to the molecular Schrödinger equation. We observe that to leading order in ϵ , the Landau-Zener formula correctly describes the probabilities for the system to remain in the original electronic level or to make a transition to the other electronic level involved in the avoided crossing. To apply the Landau-Zener formula in this case, one treats the nuclei as classical point particles to obtain a time-dependent Hamiltonian for the electrons. This leads to the study of the adiabatic limit of an effective time-dependent system with two levels isolated in its spectrum. The transition probability between the levels of such systems in the adiabatic limit is known for a variety of situations [24-26], [21], [28-31]. In particular, when the levels display an avoided crossing, the Landau-Zener formula is valid [14], [27], [22], [37], [23]. In our case, we can apply the Landau-Zener formula to the resulting time-dependent Schrödinger equation for the electrons alone.

More precisely, suppose there is a generic Type 1 or Type 2 avoided crossing at nuclear configuration x = 0. In an appropriate coordinate system, the gap between the electron energy levels is

$$2\sqrt{(b_1x_1+b_2\epsilon)^2+(c_2\epsilon)^2} + O(x^2+\epsilon^2).$$
(1.1)

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with b_1 and c_2 non-zero. Suppose that a semiclassical nuclear wave packet passes through the avoided crossing with velocity μ , whose first component is $\mu_1 \neq 0$. Then the probability of remaining in the same electronic state is

$$1 - e^{-\pi c_2^2/(b_1\mu_1)} + O(\epsilon^p), \qquad (1.2)$$

for some p > 0, and the probability of making a transition to the other electronic level involved in the avoided crossing is

$$e^{-\pi c_2^2/(b_1\mu_1)} + O(\epsilon^p). \tag{1.3}$$

This result is completely different from what one obtains for the four other types of avoided crossings, where the corresponding probabilities depend on the particular nuclear wave packet involved [19]. The reason is that in Type 1 and Type 2 avoided crossings, every part of the nuclear wave packet passes through the same size minimum gap between the eigenvalues. In Types 3, 4, 5, and 6 avoided crossings, different parts of the nuclear wave packet feel different size gaps as they pass through the avoided crossing.

Our results are also completely different from those obtained in the case of true level crossings [16]. For codimension 1 crossings, transition amplitudes produced by crossings are of order ϵ , not order 1 as in the present paper. Furthermore, the dependence of transition amplitudes on the nuclear velocity μ is exponential as opposed to algebraic in the crossing case. Also, although the problem we study in this paper is somewhat less singular than the case of true crossings, the technical details are more difficult. This stems largely from the complicated ϵ dependence of the classical mechanics in our problem. In [16], the classical mechanics had no ϵ dependence.

The Hamiltonian for a molecular system with K nuclei and N - K electrons has the form

$$H(\epsilon) = \sum_{j=1}^{K} -\frac{\epsilon^4}{2M_j} \Delta_{x_j} - \sum_{j=K+1}^{N} \frac{1}{2m_j} \Delta_{x_j} + \sum_{i< j} V_{ij}(x_i - x_j).$$
(1.4)

Here $x_j \in \mathbb{R}^l$ denotes the position of the j^{th} particle, the mass of the j^{th} nucleus is M_j (for $1 \le j \le K$), the mass of the j^{th} electron is m_j (for $K + 1 \le j \le N$), and V_{ij} is the potential between particles *i* and *j*. For convenience we assume $M_j = 1$ for $1 \le j \le K$. We set n = Kl and

let $x = (x_1, x_2, ..., x_{Kl}) \in \mathbb{R}^n$ denote the nuclear configuration vector. We decompose $H(\epsilon)$ as

$$H(\epsilon) = -\frac{\epsilon^4}{2}\Delta_x + h(x).$$
 (1.5)

This defines the electronic Hamiltonian h(x) that depends parametrically on x.

The time-dependent Schrödinger equation that we study is

$$i\epsilon^2 \frac{\partial \psi}{\partial t} = H(\epsilon)\psi,$$
 (1.6)

for t in a fixed interval. The factor of ϵ^2 on the left hand side of this equation indicates a particular choice of time scaling. Other choices could be made, but this choice is the "distinguished limit" [1] that produces the most interesting leading order solutions. With this scaling, all terms in the equation play significant roles at leading order, and the nuclear motion has a non-trivial classical limit. This is also the scaling for which the mean initial nuclear kinetic energy is held constant as ϵ tends to zero.

In this paper we are interested in the simplest types of electronic transitions that are not associated with level crossings. If the Hamiltonian has the form (1.5), then to arbitrarily high order in powers of ϵ , the solutions to (1.6) have no electronic transitions [8], [11]. There are no known rigorous results about infinite order processes in the time-dependent Born-Oppenheimer approximation. However, in real molecular systems, only a single value of ϵ is usually of interest, and h(x) may have two eigenvalues that approach one another with a minimum gap size that is of the same order of magnitude as the relevant value of ϵ . Under these circumstances, there can be significant transitions between electon energy levels. To generate useful rigorous information about these transitions without the difficulties of going beyond infinite order, we assume the electron Hamiltonian of interest can be embedded in an ϵ -dependent family that has a crossing when $\epsilon = 0$. This is in the same spirit as dealing with quantum mechanical resonances as perturbations of eigenvalues embedded in the continuous spectrum.

Thus, we study solutions to (1.6) where the Hamiltonian has the form

$$H(\epsilon) = -\frac{\epsilon^4}{2}\Delta_x + h(x, \epsilon), \qquad (1.7)$$

with the assumption that $h(x, \epsilon)$ has an Avoided Crossing according to the following definition:

DEFINITION. – Suppose $h(x, \epsilon)$ is a family of self-adjoint operators with a fixed domain \mathcal{D} in a Hilbert Space \mathcal{H} , for $x \in \Omega$ and $\epsilon \in [0, \alpha)$, where Ω is an open subset of \mathbb{R}^n . Suppose that the resolvent of $h(x, \epsilon)$ is a C^4 function of x and ϵ as an operator from \mathcal{H} to \mathcal{D} . Suppose $h(x, \epsilon)$ has two eigenvalues $E_{\mathcal{A}}(x, \epsilon)$ and $E_{\mathcal{B}}(x, \epsilon)$ that depend continuously on x and ϵ and are isolated from the rest of the spectrum of $h(x, \epsilon)$. Assume $\Gamma = \{x : E_{\mathcal{A}}(x, 0) = E_{\mathcal{B}}(x, 0)\}$ is a single point or nonempty connected proper submanifold of Ω , but that for all $x \in \Omega$, $E_{\mathcal{A}}(x, \epsilon) \neq E_{\mathcal{B}}(x, \epsilon)$ when $\epsilon > 0$. Then we say $h(x, \epsilon)$ has an Avoided Crossing on Γ .

Remark. – Realistic molecules have Coulomb potentials which give rise to electron Hamiltonians that do not satisfy the smoothness assumptions of this definition. However, one should be able to accommodate Coulomb potentials by using the regularization techniques of [10], [11], [35].

Avoided Crossings of minimal multiplicity energy levels are classified in [18], and normal forms for the electron Hamiltonian near Γ are derived. When Γ has codimension 1, there are two types. To describe these, we need some notation. Assume without loss of generality that 0 is a generic point of Γ , and decompose

$$h(x,\epsilon) = h_{\parallel}(x,\epsilon) + h_{\perp}(x,\epsilon) \tag{1.8}$$

with

$$h_{\parallel}(x,\epsilon) = h(x,\epsilon)P(x,\epsilon) \tag{1.9}$$

and

$$h_{\perp}(x,\epsilon) = h(x,\epsilon)(\mathbb{1} - P(x,\epsilon)) \tag{1.10}$$

where $P(x, \epsilon)$ is a spectral projector of $h(x, \epsilon)$ associated with $E_{\mathcal{A}}(x, \epsilon)$ and $E_{\mathcal{B}}(x, \epsilon)$. In a Type 1 Avoided Crossing, Γ has codimension 1 and the two eigenvalues each have multiplicity 1. There exists [18] an orthonormal basis $\{\psi_1(x, \epsilon), \psi_2(x, \epsilon)\}$ of $P(x, \epsilon)\mathcal{H}$, which is regular in (x, ϵ) around (0, 0). In this basis, $h_{\parallel}(x, \epsilon)$ has the form

$$h_{\parallel}(x,\epsilon) = h_1(x,\epsilon) + \widetilde{V}(x,\epsilon) = \begin{pmatrix} \beta(x,\epsilon) & \gamma(x,\epsilon) + i\delta(x,\epsilon) \\ \gamma(x,\epsilon) - i\delta(x,\epsilon) & \beta(x,\epsilon) \end{pmatrix} + \widetilde{V}(x,\epsilon) \quad (1.11)$$

where $\tilde{V}(x,\epsilon) = \frac{1}{2} \operatorname{trace}(h(x,\epsilon)P(x,\epsilon))$ is a regular function of (x,ϵ) around the origin and

$$\left.\begin{array}{l}
\left.\beta(x,\epsilon) = b_1 x_1 + b_2 \epsilon + \mathcal{O}(2) \\
\gamma(x,\epsilon) = c_2 \epsilon + \mathcal{O}(2) \\
\delta(x,\epsilon) = \mathcal{O}(2) \\
\widetilde{V}(x,\epsilon) = \mathcal{O}(0)
\end{array}\right\}$$
(1.12)

where $b_1 > 0, c_2 > 0, b_2 \in \mathbb{R}$ and

$$\mathcal{O}(m) = \mathcal{O}\left(\left(\sum_{j=0}^{n} x_j^2\right)^{m/2}\right), \quad x_0 = \epsilon.$$
 (1.13)

Type 2 Avoided Crossings are similar, except that the minimal multiplicity of eigenvalues allowed by the symmetry group is 2. Near one of these avoided crossings, one can choose an orthonormal basis $\{\psi_1(x,\epsilon), \psi_2(x,\epsilon), \psi_3(x,\epsilon), \psi_4(x,\epsilon), \}$ of $P(x,\epsilon)\mathcal{H}$, which is regular in (x,ϵ) around (0,0). In this basis, $h_{\parallel}(x,\epsilon)$ has the form

$$\begin{split} h_{\parallel}(x,\epsilon) &= h_{1}(x,\epsilon) + \widetilde{V}(x,\epsilon) = \\ \begin{pmatrix} \beta(x,\epsilon) & \gamma(x,\epsilon) + i\delta(x,\epsilon) & 0 & 0\\ \gamma(x,\epsilon) - i\delta(x,\epsilon) & -\beta(x,\epsilon) & 0 & 0\\ 0 & 0 & \beta(x,\epsilon) & \gamma(x,\epsilon) + i\delta(x,\epsilon)\\ 0 & 0 & \gamma(x,\epsilon) - i\delta(x,\epsilon) & -\beta(x,\epsilon) \\ + \widetilde{V}(x,\epsilon), & (1.14) \end{split}$$

where β , γ , δ , and \widetilde{V} are as above [18].

Throughout the paper we assume that the nuclei move transversally through Γ at the point 0. This means the classical momentum associated with their semiclassical wave packets has a non-trivial component in the x_1 direction when their classical position is passing through $0 \in \Gamma$.

Precise statements of our results require a considerable amount of notation and are presented in Theorems 3.1 and 4.1 for Type 1 and Type 2 Avoided Crossings, respectively. We have stated these theorems with the incoming state associated with the lower of the two relevant levels. The analogous results with the incoming state associated with the upper level are also true and proved in the same way, with the obvious changes. Immediate corollaries of the two theorems are that the Landau-Zener formula described above gives the correct transition probabilities in

each case. The main technique we use is matched asymptotic expansions. We use the standard time-dependent Born-Oppenheimer approximate solutions to the Schrödinger equation when the nuclei are far enough away from Γ . We match these to "inner" solutions when the system is near Γ and the standard approximation breaks down.

The paper is organized as follows: In Section 2 we discuss the ordinary differential equations whose solutions will be used to describe the semiclassical motion of the nuclei. In Section 3 we discuss semiclassical nuclear wave packets and adiabatic motion of the electrons. We then prove our main result for Type 1 Avoided Crossings, Theorem 3.1 by using matched asymptotic expansions. In Section 4 we state our main result, Theorem 4.1, for Type 2 Avoided Crossings and describe the modifications of Section 3 that are required to prove this result.

1.1. A Convenient Change of Variables

We consider the Schrödinger equation

$$i\epsilon^2 \frac{\partial}{\partial t}\psi(x,t) = -\frac{\epsilon^4}{2}\Delta\psi(x,t) + h(x,\epsilon)\psi(x,t)$$
(1.15)

In order to get rid of ϵ -dependence in the leading order of $\beta(x, \epsilon)$ in (1.11), we introduce the new variables

$$x' = b_1 x + \begin{pmatrix} b_2 \epsilon \\ 0 \\ \vdots \\ 0 \end{pmatrix}, \quad \epsilon' = c_2 \epsilon , \quad t' = b_1^2 / c_2^2 t.$$
 (1.16)

In terms of these new variables, the Schrödinger equation (1.15) for

$$\phi(x',t') = \psi(x(x',\epsilon'),t(t'))$$
(1.17)

becomes

$$i\epsilon^{\prime 2}\frac{\partial}{\partial t^{\prime}}\phi(x^{\prime},t^{\prime}) = -\frac{\epsilon^{\prime 4}}{2}\Delta\phi(x^{\prime},t^{\prime}) + \frac{c_2^4}{b_1^2}h(x(x^{\prime},\epsilon^{\prime}),\epsilon(\epsilon^{\prime}))\phi(x^{\prime},t^{\prime}) \quad (1.18)$$

in the limit $\epsilon' \rightarrow 0$, with

$$h_{\parallel}(x(x',\epsilon'),\epsilon(\epsilon')) = \begin{pmatrix} x'_1 & \epsilon' \\ \epsilon' & -x'_1 \end{pmatrix} + \mathcal{O}(2) + \widetilde{V}(x(x',\epsilon'),\epsilon(\epsilon')) \quad (1.19)$$

where $\tilde{V}(x(x', \epsilon'), \epsilon(\epsilon'))$ is regular in (x', ϵ') around (0, 0) and $\mathcal{O}(2)$ refers to x' and ϵ' . We introduce the fixed parameter $r = c_2^4/b_1^2 > 0$ and henceforth

drop the primes on the new variables. We assume that $h_1(x, \epsilon)$ has the form (1.11) with the following local behavior around x = 0 and $\epsilon = 0$:

$$\begin{cases} \beta(x,\epsilon) = rx_1 + \mathcal{O}(2) \\ \gamma(x,\epsilon) = r\epsilon + \mathcal{O}(2) \\ \delta(x,\epsilon) = \mathcal{O}(2) \\ \widetilde{V}(x,\epsilon) = \mathcal{O}(0) \end{cases}$$

$$(1.20)$$

with r > 0.

2. Ordinary Differential Equations of Semiclassical Mechanics

In Section 3.1 we introduce semiclassical wave packets for the nuclei. The leading order semiclassical motion for these wave packets is determined by the solutions to certain systems of ordinary differential equations. These involve classical mechanics, the classical action associated with a classical trajectory, and the dynamics of certain matrices that describe the position and momentum uncertainties of the wave packets. The goal of this section is to study the small ϵ behavior of these classical quantities that we need for the asymptotic matching procedure that we use in Section 3 to prove our main results.

We define

$$V^{\mathcal{A}}_{\mathcal{B}}(x,\epsilon) = \widetilde{V}(x,\epsilon) \pm \sqrt{\beta^2(x,\epsilon) + \gamma^2(x,\epsilon) + \delta^2(x,\epsilon)}$$
(2.1)

where $x \in \mathbb{R}^n$, $\epsilon > 0$. Let $a^{\mathcal{C}}(t)$ and $\eta^{\mathcal{C}}(t)$ be the solutions of the classical equations of motion

$$\left. \begin{cases} \frac{d}{dt} a^{\mathcal{C}}(t) = \eta^{\mathcal{C}}(t) \\ \frac{d}{dt} \eta^{\mathcal{C}}(t) = -\nabla V^{\mathcal{C}}(a^{\mathcal{C}}(t), \epsilon), \quad \mathcal{C} = \mathcal{A}, \quad \mathcal{B}. \end{cases} \right\}$$
(2.2)

with initial conditions

$$\left. \begin{array}{l} a^{\mathcal{C}}(0) = 0 \\ \eta^{\mathcal{C}}(0) = \eta^{0}(\epsilon), & \text{with} \\ \eta^{0}(\epsilon) = \eta^{0} + \mathcal{O}(\epsilon) , & \text{and } \eta^{0}_{1} > 0, \end{array} \right\}$$

$$(2.3)$$

where the $\mathcal{O}(\epsilon)$ term depends on whether \mathcal{C} is \mathcal{A} or \mathcal{B} . Noticing that it follows from (1.20) that $|\beta(x,\epsilon)|, |\gamma(x,\epsilon)|$ and $|\delta(x,\epsilon)|$ are $\mathcal{O}(0)$, and using estimates of the type $\beta/\sqrt{\beta^2 + \gamma^2 + \delta^2} \leq 1$, we see that

$$\|\nabla V^{\mathcal{C}}(x,\epsilon)\| = \mathcal{O}(0). \tag{2.4}$$

This last condition implies the existence and uniqueness of the solutions of (2.2).

The small ϵ perturbation theory for solutions to (2.2) and (2.3) is not quite simple because of the presence of two different time scales. However, making use of the local expressions (1.20) in the potentials $V^{\mathcal{C}}(x, \epsilon)$, we derive an asymptotic formula which holds uniformly as both t and ϵ tend to zero. An alternative, more systematic way to overcome the difficulties due to the different time scales is to use matched asymptotic expansions derived in different time regimes which agree in a non-void matching window.

2.1. Small t and ϵ Asymptotics

In order to get started, we need preliminary information about the behavior of the solution of (2.2), (2.3) when both |t| and ϵ are small.

LEMMA 2.1. – Let $a^{\mathcal{C}}(t)$ and $\eta^{\mathcal{C}}(t)$ be the solutions of (2.2) and (2.3). If ϵ and t are small enough, we have

$$\begin{cases} a^{\mathcal{C}}(t) = \eta^{0}(\epsilon)t + \mathcal{O}(t^{2}) \\ \eta^{\mathcal{C}}(t) = \eta^{0}(\epsilon) + \mathcal{O}(t) \end{cases}$$

as $t \rightarrow 0$, uniformly in ϵ .

Proof. – We mimic the proof of [16], p.82. Let us drop the index C in the notation. We want to show that there exists T > 0, such that

$$\begin{array}{l} a(t) = \eta^{0}(\epsilon)t + u(t,\epsilon) \\ b(t) = \eta^{0}(\epsilon) + v(t,\epsilon) \end{array}$$

$$(2.5)$$

where $u(t, \epsilon)/t^2$ and $v(t, \epsilon)/t$ are uniformly bounded if $\epsilon^2 + t^2 < T^2$, *i.e.*, if $(t, \epsilon) \in B_T$. We let Y_T be the Banach space of pairs of bounded, continuous in t, vector valued functions for $(t, \epsilon) \in B_T$, with the norm

$$\left\| \begin{pmatrix} u \\ v \end{pmatrix} \right\|_{Y_T} = \sup_{(t,\epsilon) \in B_T} \left\| u(t,\epsilon)/t^2 \right\| + \sup_{(t,\epsilon) \in B_T} \left\| v(t,\epsilon)/t \right\|.$$
(2.6)

For $\begin{pmatrix} u \\ v \end{pmatrix} \in Y_T$, we define

$$\mathcal{F}\begin{pmatrix} u\\v \end{pmatrix} = \begin{pmatrix} \int_0^t v(s,\epsilon)ds\\ -\int_0^t \nabla V(\eta^0(\epsilon)s + u(s,\epsilon),\epsilon)ds \end{pmatrix}.$$
 (2.7)

Due to the estimate (2.4), \mathcal{F} maps Y_T into Y_T , and furthermore, any fixed point of \mathcal{F} gives rise to a solution of (2.2), (2.3) by way of (2.5). If we show that \mathcal{F} is a strict contraction, it follows from the contraction mapping principle that the solution of (2.2), (2.3) exists, is unique and satisfies the assertion of the lemma. We show that this is the case for Tsmall enough. Let $\binom{u_1}{v_1}, \binom{u_2}{v_2} \in Y_T$ and let us estimate the norm of $\binom{u_3}{v_3} = \mathcal{F}\binom{u_1}{v_1} - \mathcal{F}\binom{u_2}{v_2}$.

$$\|u_{3}(t,\epsilon)/t^{2}\| = \frac{1}{t^{2}} \left\| \int_{0}^{t} v_{1}(s,\epsilon) - v_{2}(s,\epsilon) ds \right\|$$

$$\leq \frac{1}{t^{2}} \int_{0}^{|t|} |s| \sup_{\substack{(t,\epsilon) \in B_{T} \\ (t,\epsilon) \in B_{T}}} \|v_{1}(t,\epsilon) - v_{2}(t,\epsilon)\| ds$$

$$\leq \frac{1}{2} \left\| \begin{pmatrix} u_{1} \\ v_{1} \end{pmatrix} - \begin{pmatrix} u_{2} \\ v_{2} \end{pmatrix} \right\|_{Y_{T}}.$$
(2.9)

For the other component we use the mean value theorem

$$\nabla V(\eta^{0}(\epsilon)s + u_{1}(s,\epsilon),\epsilon) - \nabla V(\eta^{0}(\epsilon)s + u_{2}(s,\epsilon),\epsilon)$$
$$= V^{(2)}(\zeta(s,\epsilon),\epsilon)(u_{1}(s,\epsilon) - u_{2}(s,\epsilon))$$

where $V^{(2)}$ stands for the Hessian of V and

$$\zeta(s,\epsilon) = \eta^0(\epsilon)s + u_1(s,\epsilon) + \theta(s,\epsilon)(u_2(s,\epsilon) - u_1(s,\epsilon))$$
(2.11)

with $\theta(s,\epsilon) \in]0,1[$. There will appear several constants, independent of ϵ in the sequel which we shall denote generically by c. By hypothesis, $u_j(s,\epsilon) \leq cs^2, j = 1,2$, uniformly in ϵ and $\eta_1^0(\epsilon) = \eta_1^0 + \mathcal{O}(\epsilon)$ with $\eta_1^0 > 0$. Consequently, there exists a constant c, independent of ϵ , such that

$$|\zeta_1(s,\epsilon)| \ge c|s|. \tag{2.12}$$

By explicit computation we get, (omitting the arguments (x, ϵ))

$$V^{(2)} = \widetilde{V}^{(2)} \pm (|\nabla\beta\rangle\langle\nabla\beta| + |\nabla\gamma\rangle\langle\nabla\gamma| + |\nabla\delta\rangle\langle\nabla\delta|)(\beta^{2} + \gamma^{2} + \delta^{2})^{-1/2} + (\beta\beta^{(2)} + \gamma\gamma^{(2)} + \delta\delta^{(2)})(\beta^{2} + \gamma^{2} + \delta^{2})^{-1/2} - |\beta\nabla\beta + \gamma\nabla\gamma + \delta\nabla\delta\rangle\langle\beta\nabla\beta + \gamma\nabla\gamma + \delta\nabla\delta|(\beta^{2} + \gamma^{2} + \delta^{2})^{-3/2}.$$
(2.13)

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Using the behaviors (1.20) again and estimates of the type $\beta\gamma/(\beta^2 + \gamma^2 + \delta^2) \leq 1$, we get

$$\|V^{(2)}(x,\epsilon)\| \le \frac{c}{\sqrt{\beta^2(x,\epsilon) + \gamma^2(x,\epsilon) + \delta^2(x,\epsilon)}}$$
(2.14)

where c is independent of (x, ϵ) . Then,

$$\beta^{2}(x,\epsilon) + \gamma^{2}(x,\epsilon) + \delta^{2}(x,\epsilon) = r^{2}(x_{1}^{2} + \epsilon^{2}) + \mathcal{O}(3), \qquad (2.15)$$

with $\|\zeta(s,\epsilon)\| = \mathcal{O}(s)$ and (2.12) we get

$$\beta^{2}(x,\epsilon) + \gamma^{2}(x,\epsilon) + \delta^{2}(x,\epsilon)|_{x=\zeta(s,\epsilon)} = r^{2} \left(\zeta_{1}^{2}(s,\epsilon) + \epsilon^{2} \right) + \mathcal{O}(s^{3} + \epsilon^{3})$$

$$\geq c \left(\zeta_{1}^{2}(s,\epsilon) + \epsilon^{2} \right)$$

$$\geq cs^{2} \qquad (2.16)$$

if s is small enough. Hence

$$\|V^{(2)}(\zeta(s,\epsilon),\epsilon)(u_2(s,\epsilon) - u_1(s,\epsilon))\| \le c \|u_2(s,\epsilon) - u_1(s,\epsilon)\|/|s|$$
$$\le c|s| \left\| \begin{pmatrix} u_1\\v_1 \end{pmatrix} - \begin{pmatrix} u_2\\v_2 \end{pmatrix} \right\|_{Y_T}, \quad (2.17)$$

for some c independent of ϵ . Thus

$$\begin{aligned} \|v_3(t,\epsilon)\|/|t| &\leq c \frac{1}{|t|} \int_0^{|t|} |s| \left\| \begin{pmatrix} u_1 \\ v_1 \end{pmatrix} - \begin{pmatrix} u_2 \\ v_2 \end{pmatrix} \right\|_{Y_T} \\ &\leq c \frac{|t|}{2} \left\| \begin{pmatrix} u_1 \\ v_1 \end{pmatrix} - \begin{pmatrix} u_2 \\ v_2 \end{pmatrix} \right\|_{Y_T}. \end{aligned}$$
(2.18)

Choosing T small enough so that $c|t|/2 \le 1/2$, we get the result. \Box

If we replace the potential $V^{\mathcal{C}}(x,\epsilon)$ by $\widetilde{V}(x,\epsilon)$ which is regular as x and $\epsilon \rightarrow 0$, we can go further in the asymptotics, as is easily checked.

LEMMA 2.2. – Let a(t) and $\eta(t)$ be the solutions of (2.2) and (2.3) with $V^{\mathcal{C}}(x,\epsilon) \equiv \widetilde{V}(x,\epsilon)$. If ϵ and t are small enough, we have

$$\begin{cases} a(t) = \eta^{0}(\epsilon)t - \nabla \widetilde{V}(0,\epsilon)\frac{t^{2}}{2} + \mathcal{O}(t^{3}) \\ \eta(t) = \eta^{0}(\epsilon) - \nabla \widetilde{V}(0,\epsilon)t + \mathcal{O}(t^{2}) \end{cases}$$

as $t \rightarrow 0$, uniformly in ϵ .

We can now go further in the asymptotics of the classical position and velocity as both t and ϵ tend to zero.

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PROPOSITION 2.1. – Let $a^{\mathcal{C}}(t)$ and $\eta^{\mathcal{C}}(t)$ be the solutions of (2.2) and (2.3). For t and ϵ small enough, we have the asymptotics

$$a^{\frac{A}{B}}(t) = -\nabla \widetilde{V}(0,\epsilon) \frac{t^2}{2} + \eta^0(\epsilon)t + \mathcal{O}(|t|^3 + \epsilon t^2)$$

$$\mp \begin{pmatrix} 1\\0\\\vdots\\0 \end{pmatrix} \frac{r}{2\eta_1^0(\epsilon)} \left[t\sqrt{(\eta_1^0(\epsilon)t)^2 + \epsilon^2} + \frac{\epsilon^2}{\eta_1^0(\epsilon)} \ln(\eta_1^0(\epsilon)t + \sqrt{(\eta_1^0(\epsilon)t)^2 + \epsilon^2}) - \frac{\epsilon^2\ln(\epsilon)}{\eta_1^0(\epsilon)} - 2\epsilon t \right] (2.19)$$

The asymptotics for $\eta^{\mathcal{C}}(t)$ in the same regime are obtained by termwise differentiation of the above formulae up to errors $\mathcal{O}(t^2 + \epsilon |t|)$.

Proof. – By explicit computation we get, (omitting the arguments (x, ϵ))

$$-\nabla V_{\mathcal{B}}^{\mathcal{A}} = -\nabla \widetilde{V} \mp \frac{\beta \nabla \beta + \gamma \nabla \gamma + \delta \nabla \delta}{\sqrt{\beta^2 + \gamma^2 + \delta^2}}.$$
 (2.20)

Introducing the local behaviors (1.20) and replacing x by $a^{\mathcal{C}}(t)$, we make use of lemma 2.1 and $\eta^{0}(\epsilon) = \mathcal{O}(\epsilon^{0})$ (so that $\mathcal{O}(n) = \mathcal{O}(|t|^{n} + \epsilon^{n})$) to get

$$\begin{aligned} \frac{d^2}{dt^2} a^{\mathcal{C}}(t) &= -\nabla \widetilde{V}(0,\epsilon) + \mathcal{O}(t) \mp \frac{\begin{pmatrix} 1\\0\\ \vdots\\0 \end{pmatrix}} r^2 \eta_1^0(\epsilon) t + \mathcal{O}(t^2 + \epsilon^2) \\ \frac{\sqrt{(r\eta_1^0(\epsilon)t)^2 + (r\epsilon)^2 + \mathcal{O}(|t|^3 + \epsilon^3)}}{\sqrt{(r\eta_1^0(\epsilon)t)^2 + (r\epsilon)^2 + \mathcal{O}(|t|^3 + \epsilon^3)}} \\ &= -\nabla \widetilde{V}(0,\epsilon) \mp \begin{pmatrix} 1\\0\\ \vdots\\0 \end{pmatrix} \frac{r\eta_1^0(\epsilon)t}{\sqrt{((\eta_1^0(\epsilon)t)^2 + \epsilon^2)(1 + \mathcal{O}(|t| + \epsilon))}} \\ &+ \mathcal{O}(|t| + \epsilon) \\ &= -\nabla \widetilde{V}(0,\epsilon) \mp \begin{pmatrix} 1\\0\\ \vdots\\0 \end{pmatrix} \frac{r\eta_1^0(\epsilon)t}{\sqrt{(\eta_1^0(\epsilon)t)^2 + \epsilon^2}} + \mathcal{O}(|t| + \epsilon) \quad (2.21) \end{aligned}$$

We get the result by explicit integration, taking into account the initial conditions (2.3). \Box

In the sequel, we will actually need such asymptotic behaviors for matching in the time regime defined by t such that $\epsilon \rightarrow 0$, $t \rightarrow 0$, $|t|/\epsilon \rightarrow \infty$ and $t^3/\epsilon^2 \rightarrow 0$. We will refer to this regime as the matching regime.

COROLLARY 2.1. – Further expanding, we get in the matching regime $\epsilon \rightarrow 0, t \rightarrow 0, |t|/\epsilon \rightarrow \infty$ and $t^3/\epsilon^2 \rightarrow 0$

$$\begin{aligned} a^{\frac{A}{B}}(t) &= -\nabla \widetilde{V}(0,\epsilon) \frac{t^2}{2} + \eta^0(\epsilon)t \pm \frac{r}{\eta_1^0(\epsilon)} \epsilon t \begin{pmatrix} 1\\0\\\vdots\\0 \end{pmatrix} \\ &\mp \operatorname{sign}(t) r \left[\frac{t^2}{2} + \frac{\epsilon^2 \ln|t|}{2(\eta_1^0(\epsilon))^2} + \frac{\epsilon^2}{4(\eta_1^0(\epsilon))^2} (1 + 2\ln(2\eta_1^0(\epsilon))) \right. \\ &\left. - \frac{\epsilon^2 \ln \epsilon}{2(\eta_1^0(\epsilon))^2} \right] \begin{pmatrix} 1\\0\\\vdots\\0 \end{pmatrix} \\ &+ \mathcal{O}(t^3) + \mathcal{O}(\epsilon^4/t^2) \end{aligned}$$

The asymptotics for $\eta^{\mathcal{C}}(t)$ in the same regime are obtained by termwise differentiation of the above formulae up to errors $\mathcal{O}(t^2) + \mathcal{O}(\epsilon^4/t^3)$.

2.2. Classical Action Integrals

In Section 4 we construct quantum mechanical wave functions by using matched asymptotic expansions. To do so, we need the small |t| asymptotics of classical action integrals. Let

$$S^{\mathcal{C}}(t) = \int_{0}^{t} \left(\frac{\eta^{\mathcal{C}^{2}}(t')}{2} - V^{\mathcal{C}}(a^{\mathcal{C}}(t'), \epsilon) \right) dt'$$
(2.22)

$$= \int_0^t \eta^{\mathcal{C}^2}(t')dt' - \eta^{0^2}(\epsilon)\frac{t}{2} - V^{\mathcal{C}}(a^{\mathcal{C}}(0),\epsilon)t, \qquad (2.23)$$

and let S(t) be the same quantity for $V^{\mathcal{C}}(x,\epsilon) \equiv \widetilde{V}(x,\epsilon)$. From lemma 2.2 we easily deduce

LEMMA 2.3. – As $t \rightarrow 0$,

$$S(t) = \eta^{0^2}(\epsilon)\frac{t}{2} - \widetilde{V}(0,\epsilon)t - \eta^0(\epsilon)\nabla\widetilde{V}(0,\epsilon)t^2 + \mathcal{O}(t^3), \qquad (2.24)$$

uniformly in ϵ .

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From corollary 2.1, and the formula

$$V^{\mathcal{A}}_{\mathcal{B}}(0,\epsilon) = \widetilde{V}(0,\epsilon) \pm \epsilon r + \mathcal{O}(\epsilon^2)$$
(2.25)

we obtain

LEMMA 2.4. – In the regime $\epsilon \rightarrow 0$, $t \rightarrow 0$, $|t|/\epsilon \rightarrow \infty$ and $t^3/\epsilon^2 \rightarrow 0$ we have the asymptotics

$$\begin{split} S^{\overset{A}{\scriptscriptstyle \mathcal{B}}}(t) = & S^{\overset{A}{\scriptscriptstyle \mathcal{B}}}_{0}(\epsilon, \operatorname{sign}(t)) - \widetilde{V}(0, \epsilon)t + \eta^{0^{2}}(\epsilon)\frac{t}{2} - \eta^{0}(\epsilon)\nabla\widetilde{V}(0, \epsilon)t^{2} \pm r\epsilon t \\ & \mp \operatorname{sign}(t) \bigg(r\eta^{0}_{1}(\epsilon)t^{2} + \frac{r}{\eta^{0}_{1}(\epsilon)}\epsilon^{2} \ln |t| \bigg) \\ & + \mathcal{O}(t^{3}) + \mathcal{O}(\epsilon^{4}/t^{2}) + \mathcal{O}(\epsilon^{3} \ln t). \end{split}$$

2.3. Different Initial Momenta

For later purposes, we assume from now on that the solution a(t) of (2.2) with $V^{\mathcal{C}}(x,\epsilon) \equiv \widetilde{V}(x,\epsilon)$ is subject to the initial conditions

$$\begin{array}{c} a(0) = 0\\ \eta(0) = \eta^0 \end{array} \right\}$$
 (2.26)

whereas the solutions $a^{\mathcal{C}}(t)$ satisfy

$$a^{\mathcal{C}}(0) = 0 \eta^{\mathcal{C}}(0) = \eta^{0}(\epsilon) = \eta^{0} + \mathcal{O}(\epsilon).$$
(2.27)

The $\mathcal{O}(\epsilon)$ term must be included in our calculations because when the electrons makes a transition from one energy level surface to another, the nuclei must compensate by making a change in their kinetic energy in order to conserve the total energy of the whole system.

We easily get the estimates

COROLLARY 2.2. – When
$$\epsilon \rightarrow 0$$
, $t \rightarrow 0$, $t^3/\epsilon^2 \rightarrow 0$ and $|t|/\epsilon \rightarrow \infty$ we have

$$\begin{split} &\eta^{\frac{A}{B}}(t)(a(t)-a^{\frac{A}{B}}(t))\\ &=\mp r\epsilon t - (\eta^{0}(\epsilon)-\eta^{0})\eta^{0}(\epsilon)t\\ &\pm \operatorname{sign}(t) \bigg[\frac{r\eta_{1}^{0}(\epsilon)}{2}t^{2} + \frac{r}{2\eta_{1}^{0}(\epsilon)} \big(\epsilon^{2}\ln|t| + \epsilon^{2}(\ln(2\eta_{1}^{0}(\epsilon))+1/2) - \epsilon^{2}\ln\epsilon\big) \bigg]\\ &+ \mathcal{O}(t^{3}) + \mathcal{O}(\epsilon^{4}/t^{2}) + \mathcal{O}(t\epsilon^{2}\ln\epsilon) \end{split}$$

and

COROLLARY 2.3. – When
$$\epsilon \rightarrow 0$$
, $t \rightarrow 0$, $t^3/\epsilon^2 \rightarrow 0$ and $|t|/\epsilon \rightarrow \infty$ we have

$$\begin{split} S^{\mathcal{A}}_{\mathcal{B}}(t) &= S^{\mathcal{A}}_{\mathcal{B}}(\epsilon, \operatorname{sign}(t))_{0} + S(t) + (\eta^{0^{2}}(\epsilon) - \eta^{0^{2}}) \frac{t}{2} \pm r\epsilon t \\ &\mp \operatorname{sign}(t) \bigg(r\eta_{1}^{0}(\epsilon)t^{2} + \frac{r}{\eta_{1}^{0}(\epsilon)}\epsilon^{2} \ln |t| \bigg) \\ &+ \mathcal{O}(t^{3}) + \mathcal{O}(\epsilon^{4}/t^{2}) + \mathcal{O}(\epsilon^{3} \ln t). \end{split}$$

2.4. Matrices $A^{\mathcal{C}}(t)$ and $B^{\mathcal{C}}(t)$

The construction of the semiclassical wave packets that describe the nuclei requires the computation of matrices which are defined by means of classical quantities. Let $A^{\mathcal{C}}(t)$ and $B^{\mathcal{C}}(t)$ be the matrix solutions of the linear system

$$\frac{d}{dt}A^{\mathcal{C}}(t) = iB^{\mathcal{C}}(t)$$

$$\frac{d}{dt}B^{\mathcal{C}}(t) = iV^{(2)}(a^{\mathcal{C}}(t),\epsilon)A^{\mathcal{C}}(t)$$
(2.28)

where $a^{\mathcal{C}}(t)$ is the solution of (2.2) and (2.3), with initial conditions

$$A^{\mathcal{C}}(0) = A_0 \\ B^{\mathcal{C}}(0) = B_0.$$
 (2.29)

To do asymptotic matching, we need the small |t| asymptotics of $A^{\mathcal{C}}(t)$ and $B^{\mathcal{C}}(t)$. We first determine the leading order behavior of $V^{(2)}(a^{\mathcal{C}}(t),\epsilon)$ for small |t| and ϵ . From (2.13) and (1.20) it is easily seen that the Hessian matrix

$$\left(\sqrt{\beta^2 + \gamma^2 + \delta^2}\right)^{(2)} = \frac{(\delta^2 + \gamma^2)|\nabla\beta\rangle\langle\nabla\beta| + \mathcal{O}(3)}{(\beta^2 + \gamma^2 + \delta^2)^{3/2}},$$
(2.30)

where we have used the same notation as earlier. More explicitly,

$$\left(\sqrt{\beta^2(x,\epsilon) + \gamma^2(x,\epsilon) + \delta^2(x,\epsilon)}\right)^{(2)} = \frac{r\epsilon^2 P + \mathcal{O}(3)}{\left(x_1^2 + \epsilon^2 + \mathcal{O}(3)\right)^{3/2}} \quad (2.31)$$

where

$$P = \begin{pmatrix} 1 & 0 & \cdots & 0 \\ 0 & 0 & \cdots & 0 \\ \vdots & \vdots & & \vdots \\ 0 & 0 & \cdots & 0 \end{pmatrix}.$$
 (2.32)

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When x is replaced by $a^{\mathcal{C}}(t) = \eta^0(\epsilon)t + \mathcal{O}(t^2)$, the error terms $\mathcal{O}(3)$ become $\mathcal{O}((|t|^3 + \epsilon^3)$ and we obtain

$$\left(a_1^{\mathcal{C}}(t)^2 + \epsilon^2 + \mathcal{O}(3)\right)^{3/2} = \left((\eta^0(\epsilon)_1 t)^2 + \epsilon^2\right)^{3/2} (1 + \mathcal{O}(|t| + \epsilon)).$$
(2.33)

This last estimate allows us to find the leading order behavior of $V^{(2)}(a^{\mathcal{C}}(t), \epsilon)$ as $\epsilon \rightarrow 0$ and $t \rightarrow 0$:

$$V^{(2)}(a^{\frac{A}{B}}(t),\epsilon) = \pm \frac{r\epsilon^2}{\left((\eta_1^0(\epsilon)t)^2 + \epsilon^2\right)^{3/2}}P + \mathcal{O}(\epsilon^0 + t^0).$$
(2.34)

Consequently, for any positive ϵ , the equation defining $A^{\mathcal{C}}$ and $B^{\mathcal{C}}$ is regular as $t \rightarrow 0$. Moreover,

PROPOSITION 2.2. – Let $A^{\mathcal{C}}(t)$ and $B^{\mathcal{C}}(t)$ be the solutions of (2.28) and (2.29). For t and ϵ small enough, we have

$$\begin{aligned} A^{\frac{2}{B}}(t) &= A_0 + \mathcal{O}(t) \\ B^{\frac{A}{B}}(t) &= B_0 \pm irPA_0 \frac{t}{\sqrt{(\eta_1^0(\epsilon)t)^2 + \epsilon^2}} + \mathcal{O}(t), \end{aligned}$$

uniformly in ϵ .

Proof. – Equations (2.28) and (2.29) are equivalent to the linear system $\frac{d}{dt} \begin{pmatrix} A^{\mathcal{C}}(t) \\ B^{\mathcal{C}}(t) \end{pmatrix} = i \begin{pmatrix} \mathbb{O} & \mathbb{1} \\ V^{(2)}(a^{\mathcal{C}}(t), \epsilon) & \mathbb{O} \end{pmatrix} \begin{pmatrix} A^{\mathcal{C}}(t) \\ B^{\mathcal{C}}(t) \end{pmatrix}, \qquad \begin{pmatrix} A^{\mathcal{C}}(0) \\ B^{\mathcal{C}}(0) \end{pmatrix} = \begin{pmatrix} A_0 \\ B_0 \\ (2.35) \end{pmatrix},$ (2.35)

where \mathbb{O} and $\mathbb{1}$ are the $n \times n$ zero and unity matrices. Moreover, the generator is continuous and uniformly bounded for all t small enough, so that the solution exists, is unique, and satisfies the integral equation

$$\begin{pmatrix} A^{\mathcal{C}}(t) \\ B^{\mathcal{C}}(t) \end{pmatrix} = \begin{pmatrix} A_0 \\ B_0 \end{pmatrix} + i \int_0^t \begin{pmatrix} B^{\mathcal{C}}(t') \\ V^{(2)}(a^{\mathcal{C}}(t'), \epsilon) A^{\mathcal{C}}(t') \end{pmatrix} dt'.$$
(2.36)

Introducing the norm

$$||D||_{\infty} = \sup_{(t,\epsilon)\in B_T} ||D(t,\epsilon)||$$
(2.37)

on the set of matrices depending on t and ϵ , we get

$$\left\| \int_{0}^{t} B^{\mathcal{C}}(t') dt' \right\| \leq \|B^{\mathcal{C}}\|_{\infty} t, \qquad (2.28)$$

$$\int_{0}^{t} V^{(2)}(a^{\mathcal{C}}(t'), \epsilon) A^{\mathcal{C}}(t') dt' \| \leq \|A^{\mathcal{C}}\|_{\infty} \int_{0}^{t} \|V^{(2)}(a^{\mathcal{C}}(t'), \epsilon)\| dt', \qquad (2.39)$$

where, by virtue of (2.34),

$$\int_{0}^{t} \|V^{(2)}(a^{\mathcal{C}}(t'),\epsilon)\|dt' \leq c \int_{0}^{t} \left(\frac{\epsilon^{2}}{\left(t'^{2}+\epsilon^{2}\right)^{3/2}}+1\right)dt'$$
$$= c\frac{t}{\sqrt{t^{2}+\epsilon^{2}}}+ct \leq c(1+T) \equiv C \quad (2.40)$$

uniformly in ϵ . Consequently, we get the estimates

$$\|A^{\mathcal{C}} - A_0\|_{\infty} \le t \|B^{\mathcal{C}} - B_0\|_{\infty} + t \|B_0\|$$
(2.41)

$$||B^{\mathcal{C}} - B_0||_{\infty} \le C||A^{\mathcal{C}} - A_0||_{\infty} + C||A_0||$$
(2.42)

which imply

$$\|A^{\mathcal{C}} - A_0\|_{\infty} \le Ct \|A^{\mathcal{C}} - A_0\|_{\infty} + t(C\|A_0\| + \|B_0\|).$$
(2.43)

Hence, if t is so small that $1 - Ct \ge 1/2$,

$$\|A^{\mathcal{C}} - A_0\|_{\infty} \le 2t(C\|A_0\| + \|B_0\|)$$
(2.44)

and

$$A^{\mathcal{C}}(t) = A_0 + \mathcal{O}(t) \tag{2.45}$$

uniformly in ϵ . The use of this estimate and (2.34) in (2.36) yields

$$B^{A}_{\mathcal{B}}(t) = B_{0} \pm i\epsilon^{2}rPA_{0} \int_{0}^{t} \frac{dt'}{\left((\eta_{1}^{0}t')^{2} + \epsilon^{2}\right)^{3/2}} + \mathcal{O}(t) + \mathcal{O}\left(\int_{0}^{t} \frac{\epsilon^{2}t'dt'}{\left((\eta_{1}^{0}t')^{2} + \epsilon^{2}\right)^{3/2}}\right) = B_{0} \pm irPA_{0} \frac{t}{\sqrt{(\eta_{1}^{0}t)^{2} + \epsilon^{2}}} + \mathcal{O}(t) + \mathcal{O}\left(\frac{\epsilon t}{\sqrt{(\eta_{1}^{0}t)^{2} + \epsilon^{2}}}\right) = B_{0} \pm irPA_{0} \frac{t}{\sqrt{(\eta_{1}^{0}t)^{2} + \epsilon^{2}}} + \mathcal{O}(t). \quad \Box \qquad (2.46)$$

COROLLARY. – In the regime $\epsilon \rightarrow 0$, $t \rightarrow 0$ and $|t|/\epsilon \rightarrow \infty$ we have

$$B^{\mathcal{A}}_{\mathcal{B}}(t) = B_0 \mp \operatorname{sign}(t) \frac{ir}{\eta_1^0(\epsilon)} P A_0 + \mathcal{O}(t) + \mathcal{O}(\epsilon^2/t^2)$$
$$\equiv B_0^{\mathcal{A}}(\operatorname{sign}(t)) + \mathcal{O}(t) + \mathcal{O}(\epsilon^2/t^2).$$

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3. Type 1 Avoided Crossings

We now have all the ingredients required to construct an asymptotic solution to the equation

$$i\epsilon^2 \frac{\partial}{\partial t}\psi(x,t) = -\frac{\epsilon^4}{2}\Delta\psi(x,t) + h(x,\epsilon)\psi(x,t)$$
(3.1)

as $\epsilon \rightarrow 0$. Let us describe the building blocks and give their main properties.

3.1. Semiclassical Nuclear Wave Packets

The semiclassical motion of the nuclei is described by means of wave packets to be thought of as centered in phase space on the classical trajectory, and of width $\mathcal{O}(\epsilon)$. These are the same wave packets that are used in [7], [16].

Let *n* denote the dimension of the space of nuclear configurations. A multi-index $l = (l_1, l_2, \ldots, l_n)$ is an ordered *n*-tuple of non-negative integers. The order of *l* is defined to be $|l| = \sum_{j=1}^{n} l_j$, and the factorial of *l* is defined to be $l! = (l_1!)(l_2!) \cdots (l_n!)$. The symbol D^l denotes the differential operator $D^l = \frac{\partial^{ll}}{(\partial x_1)^{l_1}(\partial x_2)^{l_2}\cdots(\partial x_n)^{l_n}}$, and the symbol x^l denotes the monomial $x^l = x_1^{l_1}x_2^{l_2}\cdots x_n^{l_n}$. We denote the gradient of a function f by $f^{(1)}$ and the matrix of second partial derivatives by $f^{(2)}$. We view \mathbb{R}^n as a subset of \mathbb{C}^n , and let e_i denote the i^{th} standard basis vector in \mathbb{R}^n or \mathbb{C}^n . The inner product on \mathbb{R}^n or \mathbb{C}^n is $\langle v, w \rangle = \sum_{j=1}^n \overline{v_j}w_j$.

The semiclassical wave packets we use are products of complex Gaussians and generalizations of Hermite polynomials. The generalizations of the zeroth and first order Hermite polynomials are

$$\mathcal{H}_0(x) = 1 \tag{3.2}$$

and

$$\widetilde{\mathcal{H}}_1(x) = 2 \langle v, x \rangle, \tag{3.3}$$

where v is an arbitrary non-zero vector in \mathbb{C}^n . The generalizations of the higher order Hermite polynomials are defined recursively as follows: Let v_1, v_2, \ldots, v_m be m arbitrary non-zero vectors in \mathbb{C}^n . Then

$$\mathcal{H}_{m}(v_{1}, v_{2}, \dots, v_{m}; x)$$

$$= 2 \langle v_{m}, x \rangle \widetilde{\mathcal{H}}_{m-1}(v_{1}, v_{2}, \dots, v_{m-1}; x)$$

$$- 2 \sum_{i=1}^{m-1} \langle v_{m}, \overline{v_{i}} \rangle \widetilde{\mathcal{H}}_{m-2}(v_{1}, \dots, v_{i-1}, v_{i+1}, \dots, v_{m-1}; x).$$

One can prove [7] that these functions do not depend on the ordering of the vectors v_1, v_2, \ldots, v_m . Furthermore, if the space dimension is n = 1 and the vectors v_1, v_2, \ldots, v_m are all equal to $1 \in \mathbb{C}^1$, then $\widetilde{\mathcal{H}}_m(v_1, v_2, \ldots, v_m; x)$ is equal to the usual Hermite polynomial $H_m(x)$.

Now suppose A is a complex invertible $n \times n$ matrix. We define $|A| = [AA^*]^{1/2}$, where A^* denotes the adjoint of A. By the polar decomposition theorem, there exists a unique unitary matrix U_A , such that $A = |A|U_A$. Given a multi-index l, we define the polynomial

$$\mathcal{H}_{l}(A; x) = \widetilde{\mathcal{H}}_{|l|}(\underbrace{U_{A}e_{1}, \dots, U_{A}e_{1}}_{l_{1} \ entries}, \underbrace{U_{A}e_{2}, \dots, U_{A}e_{2}}_{l_{2} \ entries}, \underbrace{U_{A}e_{n}, \dots, U_{A}e_{n}}_{l_{n} \ entries}; x)$$
(3.4)

We can now define the semiclassical wave packets $\varphi_l(A, B, \hbar, a, \eta, x)$. In the Born-Oppenheimer approximation, the role of \hbar is played by ϵ^2 .

DEFINITION. – Let A and B be complex $n \times n$ matrices with the following properties:

A and B are invertible;
$$(3.5)$$

$$BA^{-1}$$
 is symmetric ([real symmetric] + i[real symmetric]); (3.6)

$$Re BA^{-1} = \frac{1}{2}[(BA^{-1}) + (BA^{-1})^*] \text{ is strictly positive definite;} (3.7)$$

$$(Re \ BA^{-1})^{-1} = AA^*. ag{3.8}$$

Let $a \in \mathbb{R}^n$, $\eta \in \mathbb{R}^n$, and $\hbar > 0$. Then for each multi-index l we define

$$\begin{aligned} \varphi_l(A, B, \hbar, a, \eta, x) \\ &= 2^{-|l|/2} (l!)^{-1/2} \pi^{-n/4} \hbar^{-n/4} \left[\det A \right]^{-1/2} \\ &\cdot \mathcal{H}_l(A; \hbar^{-1/2} |A|^{-1} (x-a)) \\ &\cdot \exp \left\{ - \langle (x-a), BA^{-1} (x-a) \rangle / 2\hbar + i \langle \eta, (x-a) \rangle / \hbar \right\}. \end{aligned}$$

The choice of the branch of the square root of $[\det A]^{-1}$ in this definition depends on the context, and is determined by initial conditions and continuity in time.

Remarks. – 1. We never use the functions $\varphi_l(A, B, \hbar, a, \eta, x)$ unless conditions (3.5)-(3.8) are satisfied.

2. Condition (3.8) is equivalent to the more symmetrical condition

$$A^*B + B^*A = 21. (3.9)$$

3. For fixed A, B, \hbar , a, and η , the functions $\varphi_l(A, B, \hbar, a, \eta, x)$ form an orthonormal basis of $L^2(\mathbb{R}^n)$. The proof can be found in [7].

4. Generically U_A and U_B are complex unitary matrices, and A and B are Hermitian. In the special cases when they all happen to be real, the functions $\varphi_l(A, B, \hbar, a, \eta, x)$ are simply rotated and dilated eigenfunctions of the n dimensional harmonic oscillator.

5. The utility of the functions $\varphi_l(A, B, \hbar, a, \eta, x)$ stems from their remarkably beautiful behavior under Fourier transforms. If we define the scaled Fourier transform to be

$$[\mathcal{F}_{\hbar}\Psi](\xi) = (2\pi\hbar)^{-n/2} \int_{\mathbb{R}^n} \Psi(x) e^{-i\langle\xi, x\rangle/\hbar} dx, \qquad (3.10)$$

then

$$[\mathcal{F}_{\hbar} \varphi_{l}(A, B, \hbar, a, \eta, \cdot)](\xi) = (-i)^{|l|} e^{-i\langle \eta, a \rangle/\hbar} \varphi_{l}(B, A, \hbar, \eta, -a, \xi).$$
(3.11)

The only proof we know of this formula involves a messy induction on |l|. The details may be found in [7].

6. The functions $\varphi_l(A, B, \hbar, a, \eta, x)$ separate the position and momentum uncertainties from one another. For any given l, the position uncertainty depends only on |A| and the momentum uncertainty depends only on |B|. For example, when the space dimension is n = 1, the position and momentum uncertainties of $\varphi_l(A, B, \hbar, a, \eta, x)$ are given by $[(l + \frac{1}{2})\hbar]^{1/2}|A|$ and $[(l + \frac{1}{2})\hbar]^{1/2}|B|$, respectively.

7. For technical reasons related to Remark 6, it is crucial that the matrix B only appear in the exponent in the definition of $\varphi_l(A, B, \hbar, a, \eta, x)$, and not in the polynomial. By Remark 5, the matrix A only occurs in the exponent of the Fourier transform. This turns out to be technically crucial, also. The fulfillment of these technical requirements makes the Fourier transform formula in Remark 5 all the more amazing.

The formulas for the functions $\varphi_l(A, B, \hbar, a, \eta, x)$ are rather complicated, but the leading order semiclassical propagation of these wave packets is very simple. Under mild hypotheses (*e.g.*, $V \in C^3$ and bounded below), the Schrödinger equation

$$i\hbar \frac{\partial \Psi}{\partial t} = -\frac{\hbar^2}{2} \frac{\partial^2 \Psi}{\partial x^2} + V(x) \Psi$$
 (3.12)

has an approximate solution of the form

$$e^{iS(t)/\hbar} \varphi_l(A(t), B(t), \hbar, a(t), \eta(t), x) + O(\hbar^{1/2}).$$
 (3.13)

Here $O(\hbar^{1/2})$ means that the exact solution and the approximate solution agree up to an error whose norm is bounded by an *l*-dependent constant times $\hbar^{1/2}$ for t in a fixed bounded interval [-T, T]. The vectors a(t) and $\eta(t)$ satisfy the classical equations of motion

$$\frac{\partial a}{\partial t}(t) = \eta(t),$$
(3.14)
$$\frac{\partial \eta}{\partial t}(t) = -V^{(1)}(a(t)).$$

The function S(t) is the classical action integral associated with the classical path,

$$S(t) = \int_{-T}^{t} \left(\frac{(\eta(s))^2}{2} - V(a(s)) \right) ds.$$
 (3.16)

The matrices A(t) and B(t) satisfy

$$\frac{\partial A}{\partial t}(t) = iB(t), \qquad (3.17)$$

$$\frac{\partial B}{\partial t}(t) = iV^{(2)}(a(t))A(t).$$
(3.18)

If A(-T) and B(-T) satisfy conditions (3.5)-(3.8), then so do A(t) and B(t) for each t. The proofs of the claims we have made about the $\varphi_l(A, B, \hbar, a, \eta, x)$ and other properties of the quantities introduced can be found in [7].

Consider φ_l defined in (3.9). Since the vectors and matrices a, η, A and B will be replaced in these functions by $a^{\mathcal{C}}(t), \eta^{\mathcal{C}}(t), A^{\mathcal{C}}(t)$ and $B^{\mathcal{C}}(t)$ solutions of (2.2) and (2.28), we need to control the changes of φ_l induced by changes in a, η, A and B.

LEMMA 3.1. – We have, in the $L^2(\mathbb{R}^n)$ sense,

$$\begin{aligned} \varphi_l(A, B, \epsilon^2, a, 0, x) = \varphi_l(A_0, B_0, \epsilon^2, a_0, 0, x) \\ &+ \mathcal{O}(||A - A_0|| + ||B - B_0|| + ||a - a_0||/\epsilon). \end{aligned}$$

Proof. - Let us introduce the temporary notation

$$\psi(A, B, y) = \varphi_l(A, B, 1, 0, 0, y) \tag{3.19}$$

so that

$$\varphi_l(A, B, \epsilon^2, a, 0, x) = \epsilon^{-n/2} \psi(A, B, (x-a)/\epsilon).$$
(3.20)

Consider

$$\begin{split} \psi(A + A', B + B', (x - a - a')/\epsilon) &- \psi(A, B, (x - a)/\epsilon) \\ &= \psi(A + A', B + B', (x - a - a')/\epsilon) - \psi(A + A', B, (x - a - a')/\epsilon) \\ &+ \psi(A + A', B, (x - a - a')/\epsilon) - \psi(A, B, (x - a - a')/\epsilon) \\ &+ \psi(A, B, (x - a - a')/\epsilon) - \psi(A, B, (x - a)/\epsilon). \end{split}$$
(3.21)

We separately estimate each of the three differences on the right hand side of this equation.

We compute the Fréchet differential,

$$D_B\psi(A, B, y)(B') = -\frac{1}{2}\psi(A, B, y)\langle y|B'A^{-1}y\rangle, \qquad (3.22)$$

so that by the mean value theorem,

$$\begin{split} |\psi(A+A',B+B',(x-a-a')/\epsilon) - \psi(A+A',B,(x-a-a')/\epsilon)| \\ &\leq \sup_{t\in[0,1]} |\psi(A+A',B+tB',(x-a-a')/\epsilon)| \\ &\frac{|x-a-a'|^2}{\epsilon^2} \|(A+A')^{-1}\| \|B'\|. \end{split}$$

But, since by assumption $\operatorname{Re} BA^{-1}$ is strictly positive definite, and

$$h(A, |A|^{-1}y) = \frac{1}{2^{|l|/2} l!^{1/2} \pi^{n/4} (\det A)^{1/2}} \mathcal{H}_l(A, |A|^{-1}y)$$
(3.24)

is a polynomial of degree |l| in its second variable, we have, for ||A'||, ||B'|| small enough, and for any positive p,

$$\epsilon^{-n} \int_{\mathbb{R}^{n}} \left| \exp\left(-\langle (x - a - a') | (B + B') (A + A')^{-1} (x - a - a') \rangle / (2\epsilon^{2}) \right) \right|^{2} \\ \times \left| (x - a - a') / \epsilon \right|^{p} dx \\ = \int_{\mathbb{R}^{n}} \left| \exp\left(-\langle y | (B + B') (A + A')^{-1} y \rangle / 2 \right) \right|^{2} |y|^{p} dy \\ \leq n_{p} (A, B), \tag{3.25}$$

uniformly in A'. B', a and a'. Thus, for some constant depending only on A and B,

$$\epsilon^{-n/2} \| \psi(A + A', B + B', (\cdot - a - a')/\epsilon) - \psi(A + A', B, (\cdot - a - a')/\epsilon) \|_{L^2(\mathbb{R}^n)} \leq C(A, B) \| B' \|.$$
(3.26)

This estimate together with the Plancherel identity and (3.10) yield

$$\begin{aligned} \epsilon^{-n/2} \|\psi(A+A', B, (\cdot - a - a')/\epsilon) - \psi(A, B, (\cdot - a - a')/\epsilon)\|_{L^{2}(\mathbb{R}^{n})} \\ &\leq \epsilon^{-n/2} \|[\mathcal{F}_{\hbar}\Psi](A+A', B, (\cdot - a - a')/\epsilon) \\ &- [\mathcal{F}_{\hbar}\Psi](A, B, (\cdot - a - a')/\epsilon)\|_{L^{2}(\mathbb{R}^{n})} \\ &= \epsilon^{-n/2} \|\psi(B, A + A', (\cdot + a + a')/\epsilon) - \psi(B, A, (\cdot + a + a')/\epsilon)\|_{L^{2}(\mathbb{R}^{n})} \\ &\leq C(B, A) \|A'\|. \end{aligned}$$
(3.27)

Finally, using the notation given in (3.24),

...

$$|\psi(A, B, y) - \psi(A, B, y + y')| \le \sup_{t \in [0,1]} |D_y \psi(A, B, y + ty')(y')| \quad (3.28)$$

where

$$D_{y}\psi(A, B, y)(y') = -\frac{1}{2}\psi(A, B, y)(\langle y|BA^{-1}y'\rangle + \langle y'|BA^{-1}y\rangle) + \exp\left(-\langle y|BA^{-1}y\rangle/2\right)D_{2}h(A, |A|^{-1}y)(|A|^{-1}y').$$
(3.29)

Thus we can write

$$\begin{aligned} \epsilon^{-n/2} |\psi(A, B, (x - a - a')/\epsilon) - \psi(A, B, (x - a)/\epsilon)| \\ &\leq \epsilon^{-n/2} \sup_{t \in [0,1]} |\exp\left(-\langle (x - a - ta')/\epsilon | BA^{-1}(x - a - ta')/\epsilon \rangle/2\right) \\ &\qquad \times P(A, |A|^{-1}(x - a - ta')/\epsilon) |||a'||/\epsilon \\ &= \epsilon^{-n/2} |\exp\left(-\langle (x - a - t_0a')/\epsilon | BA^{-1}(x - a - t_0a')/\epsilon \rangle/2\right) \\ &\qquad \times P(A, |A|^{-1}(x - a - t_0a')/\epsilon) |||a'||/\epsilon \end{aligned}$$
(3.30)

where P(A, y) is a polynomial of order |l| in y and $[0, 1] \ni t_0 = t_0(A, B, a, a', \epsilon)$ is the point where the supremum is reached. Hence, using (3.25) again, we deduce

$$\epsilon^{-n/2} \|\psi(A, B, (x-a-a')/\epsilon) - \psi(A, B, (x-a)/\epsilon)\|_{L^2(\mathbb{R}^n)} \le C'(A, B) \|a'\|/\epsilon$$
(3.31)

for some constant C'(A, B) uniform in a, a' and t_0 . \Box

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3.2 The Choice of Eigenvectors

In this section we present a procedure for choosing the electronic eigenvectors and their phases. Although the electronic hamiltonian is independent of time, it is convenient, since we deal with the time dependent Schrödinger equation, to choose specific time dependent electronic eigenvectors. Indeed, the electrons follow the motion of the nuclei in an adiabatic way, so the suitable instantaneous electronic eigenvectors must satisfy some parallel transport condition to take into account the geometric phases arising in this situation. These eigenvectors thus depend on the classical trajectories. Since they may become singular when the corresponding eigenvalues are degenerate, or almost degenerate, we shall define them for t in the outer regime. We shall have two sets of eigenvectors, denoted by $\Phi_c^{\pm}(x, t, \epsilon)$, where the label \pm , refers to positive and negative times.

Let $\eta^{\mathcal{C}}(t)$ be the momentum solution of the classical equations of motion (2.2) and (2.3). The normalized eigenvectors $\Phi_{\mathcal{C}}^{\pm}(x, t, \epsilon)$ are the solutions of

$$\langle \Phi_{\mathcal{C}}^{\pm}(x,t,\epsilon) | \left(\partial/\partial t + \eta^{\mathcal{C}}(t) \nabla \right) \Phi_{\mathcal{C}}^{\pm}(x,t,\epsilon) \rangle \equiv 0$$
 (3.32)

for C = A, B and $t \ge 0$. Since the eigenvalues $E_A(x, \epsilon)$ and $E_B(x, \epsilon)$ are non-degenerate for any time t, t small enough, such vectors exist, are unique up to an overall time independent phase factors and are eigenvectors of $h_1(x, \epsilon)$ associated with $E_C(x, \epsilon)$ for any time, see [16]. Let us make the construction of these eigenvectors more specific and give their small t and small ϵ asymptotics. We define the angles $\varphi(x, \epsilon)$ and $\theta(x, \epsilon)$ by

$$\beta(x,\epsilon) = \sqrt{\beta^2(x,\epsilon) + \gamma^2(x,\epsilon) + \delta^2(x,\epsilon)} \cos(\theta(x,\epsilon))$$
(3.33)

$$\gamma(x,\epsilon) = \sqrt{\beta^2(x,\epsilon) + \gamma^2(x,\epsilon) + \delta^2(x,\epsilon)} \sin(\theta(x,\epsilon)) \cos(\varphi(x,\epsilon)) \quad (3.34)$$

$$\delta(x,\epsilon) = \sqrt{\beta^2(x,\epsilon) + \gamma^2(x,\epsilon) + \delta^2(x,\epsilon)\sin(\theta(x,\epsilon))\sin(\varphi(x,\epsilon))} \quad (3.35)$$

and construct static eigenvectors. Let

$$\Phi_{\mathcal{A}}^{-}(x,\epsilon) = e^{i\varphi(x,\epsilon)}\cos(\theta(x,\epsilon)/2)\psi_{1}(x,\epsilon) + \sin(\theta(x,\epsilon)/2)\psi_{2}(x,\epsilon) \quad (3.36)$$

$$\Phi_{\mathcal{B}}^{-}(x,\epsilon) = e^{-i\varphi(x,\epsilon)}\cos(\theta(x,\epsilon)/2)\psi_{2}(x,\epsilon) - \sin(\theta(x,\epsilon)/2)\psi_{1}(x,\epsilon) \quad (3.37)$$

be the eigenvectors of $h_1(x, \epsilon)$ associated with $E_{\mathcal{C}}(x, \epsilon)$, $\mathcal{C} = \mathcal{A}.\mathcal{B}$ for $\pi/2 < \theta(x, \epsilon) \leq \pi$ and

$$\Phi_{\mathcal{A}}^{+}(x,\epsilon) = \cos(\theta(x,\epsilon)/2)\psi_{1}(x,\epsilon) + e^{-i\varphi(x,\epsilon)}\sin(\theta(x,\epsilon)/2)\psi_{2}(x,\epsilon) \quad (3.38)$$

$$\Phi_{\mathcal{B}}^{+}(x,\epsilon) = \cos(\theta(x,\epsilon)/2)\psi_{2}(x,\epsilon) - e^{i\varphi(x,\epsilon)}\sin(\theta(x,\epsilon)/2)\psi_{1}(x,\epsilon) \quad (3.39)$$

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be the eigenvectors of $h_1(x,\epsilon)$ for $0 \le \theta(x,\epsilon) < \pi/2$. The solutions of (3.32) are of the form

$$\Phi_{\mathcal{C}}^{\pm}(x,t,\epsilon) = \Phi_{\mathcal{C}}^{\pm}(x,\epsilon)e^{i\lambda_{\mathcal{C}}^{\pm}(x,t,\epsilon)}, \qquad \begin{cases} t > 0\\ t < 0 \end{cases}, \tag{3.40}$$

where $\lambda_{\mathcal{C}}^{\pm}(x,t,\epsilon)$ is a real valued function satisfying the equation

$$i\frac{\partial}{\partial t}\lambda_{\mathcal{C}}^{\pm}(x,t,\epsilon) + i\eta^{\mathcal{C}}(t)\nabla\lambda_{\mathcal{C}}^{\pm}(x,t,\epsilon) + \langle\Phi_{\mathcal{C}}(x,\epsilon)|\eta^{\mathcal{C}}(t)\nabla\Phi_{\mathcal{C}}(x,\epsilon)\rangle = 0.$$
(3.41)

LEMMA 3.2. – Let $\eta^{\mathcal{C}}(t)$ be the momentum solution of the classical equations of motion (2.2) and (2.3). There exist eigenvectors $\Phi_{\mathcal{C}}^{\pm}(x,t,\epsilon)$, $\mathcal{C} = \mathcal{A}, \mathcal{B}$, such that

$$\langle \Phi^{\pm}_{\mathcal{C}}(x,t,\epsilon) | (\partial/\partial t + \eta^{\mathcal{C}}(t)\nabla) \Phi^{\pm}_{\mathcal{C}}(x,t,\epsilon) \rangle \equiv 0.$$

If $||x|| = \mathcal{O}(\epsilon^{\kappa})$, $x_1 \stackrel{>}{_{<}} 0$, with $2/3 < \kappa < 1$,

$$\left|\Phi^{\pm}_{\mathcal{C}}(x,t,\epsilon) - \Phi^{\pm}_{\mathcal{C}}(x,\epsilon)\right| = \mathcal{O}(t/\epsilon^{1-\kappa}).$$

Proof. – We give a proof for $\Phi_{\mathcal{A}}^+(x, t, \epsilon)$ only; the other cases are similar. In order to simplify the notation, we drop the indices \mathcal{A} and +. Thus we need determine $\lambda(x, t, \epsilon)$ such that

$$\Phi(x,t,\epsilon) = e^{i\lambda(x,t,\epsilon)}\Phi(x,\epsilon), \qquad (3.42)$$

where

$$\Phi(x,\epsilon) = \cos(\theta(x,\epsilon)/2)\psi_1(x,\epsilon) + e^{-i\varphi(x,\epsilon)}\sin(\theta(x,\epsilon)/2)\psi_2(x,\epsilon).$$
(3.43)

Also, let $a(t) \equiv a^{\mathcal{A}}(t)$ and $\eta(t) \equiv \eta^{\mathcal{A}}(t)$ for t > 0. We introduce the new variable

$$\omega \equiv x - a(t) \tag{3.44}$$

and the notation

$$\lambda(\omega, t, \epsilon) \equiv \lambda(\omega + a(t), t, \epsilon), \qquad (3.45)$$

$$\Phi(\omega, t, \epsilon) \equiv \Phi(\omega + a(t), \epsilon). \tag{3.46}$$

In terms of these new variables, equation (3.41) for $\tilde{\lambda}$ reads

$$i\frac{\partial}{\partial t}\widetilde{\lambda}(\omega,t,\epsilon) = -\langle \widetilde{\Phi}(\omega,t,\epsilon) | \frac{\partial}{\partial t}\widetilde{\Phi}(\omega,t,\epsilon) \rangle$$
(3.47)

with

$$\frac{\partial}{\partial t}\widetilde{\Phi}(\omega, t, \epsilon) = \eta(t)\nabla\Phi(\omega + a(t), \epsilon).$$
(3.48)

We compute, dropping the arguments,

$$\nabla \Phi = \left(-\sin(\theta/2)\psi_1 + e^{-i\varphi}\cos(\theta/2)\psi_2\right)\nabla\theta/2 + \cos(\theta/2)\nabla\psi_1 + e^{-i\varphi}\sin(\theta/2)\nabla\psi_2 - i\nabla\varphi e^{-i\varphi}\sin(\theta/2)\psi_2.$$
(3.49)

Since the ψ_j , j = 1, 2, are orthonormal, we get

$$\left\langle \tilde{\Phi} | \frac{\partial}{\partial t} \tilde{\Phi} \right\rangle = -i\eta \nabla \varphi \sin^2(\theta/2) + \cos^2(\theta/2) \langle \psi_1 | \eta \nabla \psi_1 \rangle + \sin^2(\theta/2) \langle \psi_2 | \eta \nabla \psi_2 \rangle + \sin(\theta/2) \cos(\theta/2) \left(e^{-i\varphi} \langle \psi_1 | \eta \nabla \psi_2 \rangle + e^{i\varphi} \langle \psi_2 | \eta \nabla \psi_1 \rangle \right).$$

$$(3.50)$$

As $\eta, \psi_j, \nabla \psi_j$ are all $\mathcal{O}(0)$, we can write

$$\widetilde{\lambda}(\omega, t, \epsilon) = \int^{t} \left(\eta(t') \nabla \varphi \sin^{2}(\theta(\omega + a(t'), \epsilon)/2) + \mathcal{O}(0) \right) dt' + \widetilde{\lambda}_{0}(\omega).$$
(3.51)

It remains to estimate $\nabla \varphi$.

$$\varphi(x,\epsilon) = \operatorname{atan}\left(\delta(x,\epsilon)/\gamma(x,\epsilon)\right)$$
 (3.52)

provided $\gamma(x,\epsilon)$ is different from zero. Hence, using (1.20) we get

$$\nabla\varphi(x,\epsilon) = \frac{\gamma(x,\epsilon)\nabla\delta(x,\epsilon) - \delta(x,\epsilon)\nabla\gamma(x,\epsilon)}{\gamma^2(x,\epsilon) + \delta^2(x,\epsilon)} = \frac{\epsilon\mathcal{O}(1) + \mathcal{O}(2)\mathcal{O}(1)}{r^2\epsilon^2 + \mathcal{O}(3)}$$
$$= \frac{\mathcal{O}(1)/\epsilon + \mathcal{O}(3)/\epsilon^2}{1 + \mathcal{O}(3)/\epsilon^2}.$$
(3.53)

The above computations are legitimate for

$$\|x\| = \mathcal{O}(\epsilon^{\kappa}), \qquad 2/3 < \kappa < 1, \tag{3.54}$$

which implies

$$\gamma(x,\epsilon) = r\epsilon + \mathcal{O}(2)$$

= $r\epsilon(1 + \mathcal{O}(\epsilon + ||x|| + ||x||^2/\epsilon))$
= $r\epsilon(1 + \mathcal{O}(\epsilon^{2\kappa-1}))$
 $\geq c\epsilon,$ (3.55)

for some uniform constant c, and

$$\mathcal{O}(1)/\epsilon = \mathcal{O}(1 + ||x||/\epsilon)$$

= $\mathcal{O}(1/\epsilon^{1-\kappa}),$ (3.56)

$$\mathcal{O}(3)/\epsilon^{2} = \mathcal{O}(\epsilon + ||x|| + ||x||^{2}/\epsilon + ||x||^{3}/\epsilon^{2})$$

 $\to 0$
(3.57)

as $\epsilon \rightarrow 0$. Hence, if $\|\omega + a(t)\| = \mathcal{O}(\epsilon^{\kappa})$,

$$\widetilde{\lambda}(\omega, t, \epsilon) = \widetilde{\lambda}_0(\omega) + \mathcal{O}(t/\epsilon^{1-\kappa}), \qquad (3.58)$$

uniformly in t and ω . Taking the integration constant $\widetilde{\lambda}_0(\omega) \equiv 0$, we obtain

$$\lambda(x,t,\epsilon) = \mathcal{O}(t/\epsilon^{1-\kappa}). \tag{3.59}$$

We get the desired result by means of (3.42) and the fact that $\Phi(x, \epsilon) = O(0)$. \Box

The nuclear wave function is localized around the classical trajectory in the semiclassical regime. Thus, in the outer temporal region, because of the genericity condition $\eta_1^0 > 0$, the major part of the nuclear wave function will be supported away from the neighborhood where the levels almost cross. Hence we can introduce a cutoff function which does not significantly alter the solution and forces the support of the wave function to be away from this neighborhood. We choose the support of the cutoff so that the above lemma applies. Let F be a C^{∞} cutoff function

$$F: \mathbb{R}^+ \mapsto \mathbb{R}, \tag{3.60}$$

such that

$$F(r) = 1 \quad 0 \le r \le 1 F(r) = 0 \quad r \ge 2$$
 (3.61)

The wave functions we construct below in the outer regime will be multiplied by the regulating factor

$$F(\|x - a^{\mathcal{C}}(t)\|/\epsilon^{1-\delta'})$$
(3.62)

where $\delta' < \xi$, for C = A.B (see below).

Remark. – If $\xi < 2/3 < \kappa < 1 - \xi$, the correction terms in the above lemma tend to zero and the set $||x|| = \mathcal{O}(\epsilon^{\kappa})$ includes the set $||x - a^{\mathcal{C}}(t)|| = \mathcal{O}(\epsilon^{1-\delta'})$.

On the support of F,

$$x = \eta^0(\epsilon)t + \mathcal{O}(\epsilon^{1-\delta'} + t^2)$$
(3.63)

and since $\eta_1^0(\epsilon) = \eta_1^0 + \mathcal{O}(\epsilon)$, where $\eta_1^0 > 0$, we deduce

$$|x_1| > c|t| \tag{3.64}$$

uniformly in ϵ . Using this property we have

LEMMA 3.3. – For x in the support of $F(||x - a^{\mathcal{C}}(t)||/\epsilon^{1-\delta'})$, $\mathcal{C} = \mathcal{A}, \mathcal{B}$, we have for t < 0,

$$\Phi^-_{\mathcal{A}}(x,\epsilon) = \psi_2(x,\epsilon) + \mathcal{O}(\epsilon/t^{1/2}), \ \Phi^-_{\mathcal{B}}(x,\epsilon) = -\psi_1(x,\epsilon) + \mathcal{O}(\epsilon/t^{1/2}),$$

and for t > 0,

$$\Phi^+_{\mathcal{A}}(x,\epsilon) = \psi_1(x,\epsilon) + \mathcal{O}(\epsilon/t^{1/2}),$$

$$\Phi^+_{\mathcal{B}}(x,\epsilon) = \psi_2(x,\epsilon) + \mathcal{O}(\epsilon/t^{1/2}).$$

Proof. – It follows from (1.20), (3.63) and (3.64) that for $t \stackrel{>}{_{<}} 0$,

$$\cos(\theta(x,\epsilon)) = \frac{x_1 + \mathcal{O}(2)}{\sqrt{x_1^2 + \epsilon^2 + \mathcal{O}(3)}} = \frac{\pm 1 + \mathcal{O}(t)}{(1 + \mathcal{O}(\epsilon^2/t^2) + \mathcal{O}(t))^{1/2}} = \pm 1 + \mathcal{O}(\epsilon^2/t).$$
(3.65)

Similarly,

$$\tan(\varphi(x,\epsilon)) = \frac{\mathcal{O}(2)}{\epsilon(1+\mathcal{O}(2))} = \frac{\mathcal{O}(t^2/\epsilon)}{1+\mathcal{O}(t^2/\epsilon)} = \mathcal{O}(t^2/\epsilon), \qquad (3.66)$$

hence

$$\varphi(x,\epsilon) = \mathcal{O}(t^2/\epsilon). \tag{3.67}$$

The result then follows from the identity

$$\sin^{2}(\theta(x,\epsilon)/2) = \frac{1 - \cos(\theta(x,\epsilon))}{2} = \frac{1 \pm 1}{2} + \mathcal{O}(t^{2}/\epsilon)$$
(3.68)

and $\|\psi_j\| = 1$, j = 1, 2. \Box Vol. 68, n° 1-1998. We can describe the solution of the Schrödinger equation in the different regimes introduced.

3.3. The Incoming Outer Solution

We assume that at the initial time -T, the wave function is given by a semiclassical nuclear wave packet times an electronic function associated with the \mathcal{B} level. We consider the associated classical quantities determined by the following initial conditions at t = 0:

$$\begin{array}{l}
 a^{\mathcal{B}}(0) = 0 \\
\eta^{\mathcal{B}}(0) = \eta^{0} \\
A^{\mathcal{B}}(0) = A_{0} \\
B^{\mathcal{B}}(0) = B_{0} \\
S^{\mathcal{B}}(0) = 0,
\end{array}$$
(3.69)

where A_0 and B_0 satisfy the hypotheses in the definition of the functions φ_l . Away from the avoided crossing of the electronic levels, the solution of the Schrödinger equation is well approximated by standard time-dependent Born-Oppenheimer wave packets. Close to the avoided crossing these standard wave packets fail to approximate the solution. The following lemma tells us how close to the avoided crossing time these standard wave packets can be used as approximations.

LEMMA 3.4. – In the incoming outer time regime, $-T \leq t \leq -\epsilon^{1-\xi}$, there is an approximation $\psi_{IO}(x,t)$ of a solution of the Schrödinger equation $\psi(x,t)$ of the form

$$\psi_{IO}(x,t) = F(||x-a^{\mathcal{B}}(t)||/\epsilon^{1-\delta'}) \exp\left(i\frac{S^{\mathcal{B}}(t)}{\epsilon^2} + i\frac{\eta^{\mathcal{B}}(t)(x-a^{\mathcal{B}}(t))}{\epsilon^2}\right) \times \varphi_l(A^{\mathcal{B}}(t), B^{\mathcal{B}}(t), \epsilon^2, a^{\mathcal{B}}(t), 0, x)\Phi_{\mathcal{B}}^-(x, t, \epsilon)$$
(3.70)

such that

$$\psi(x,t) = \psi_{IO}(x,t) + R(x,t,\epsilon) \tag{3.71}$$

where

$$R(x,t,\epsilon) = \mathcal{O}(\epsilon^{\xi}) \tag{3.72}$$

in the $L^2(\mathbb{R}^n)$ sense, as $\epsilon \rightarrow 0$.

Proof. – The proof of this Lemma is very similar to that of Lemma 6.4 of [16]. For $-T \leq t \leq -\epsilon^{1-\xi}$ and x in the support of the cut off

function $F(||x - a^{\mathcal{B}}(t)||/\epsilon^{1-\delta'})$, $\frac{1}{E_{\mathcal{A}}(x) - E_{\mathcal{B}}(x)}$ is bounded by a multiple of $\epsilon^{-1+\xi}$. Furthermore, for such x and t, a rather tedious calculation shows that $(\frac{\partial}{\partial t} + \eta^{\mathcal{B}} \cdot \nabla_x) \Phi_{\mathcal{B}}^-(x, t, \epsilon)$ is bounded by a multiple of $\epsilon^{-1+\xi}$. Thus, up to errors on the order of $\epsilon^{2\xi}$, $\psi_{IO}(x, t)$ is equal to

$$\Psi(x,t) = F(||x-a^{\mathcal{B}}(t)||/\epsilon^{1-\delta'}) \exp\left(i\frac{S^{\mathcal{B}}(t)}{\epsilon^2} + i\frac{\eta^{\mathcal{B}}(t)(x-a^{\mathcal{B}}(t))}{\epsilon^2}\right) \times \left(\Psi_0(x,t) + \epsilon^2 \Psi_2^{\perp}(x,t)\right)$$
(3.73)

where

$$\Psi_0(x,t) = \varphi_l(A^{\mathcal{B}}(t), B^{\mathcal{B}}(t), \epsilon^2, a^{\mathcal{B}}(t), 0, x) \Phi_{\mathcal{B}}^-(x, t, \epsilon)$$
(3.74)

and

$$\Psi_{2}^{\perp}(x,t) = i\varphi_{l}(A^{\mathcal{B}}(t), B^{\mathcal{B}}(t), \epsilon^{2}, a^{\mathcal{B}}(t), 0, x)$$

$$\times \left(\Phi_{\mathcal{A}}^{-}(x,t,\epsilon) \frac{1}{E_{\mathcal{A}}(x) - E_{\mathcal{B}}(x)} \langle \Phi_{\mathcal{A}}^{-}(x,t,\epsilon), \left(\frac{\partial}{\partial t} + \eta^{\mathcal{B}} \cdot \nabla_{x}\right) \Phi_{\mathcal{B}}^{-}(x,t,\epsilon) \rangle + r_{\mathcal{B}\mathcal{A}}(x,\epsilon) P_{\perp}^{\mathcal{A}\mathcal{B}}(x,\epsilon) \left(\frac{\partial}{\partial t} + \eta^{\mathcal{B}} \cdot \nabla_{x}\right) \Phi_{\mathcal{B}}^{-}(x,t,\epsilon) \right), \qquad (3.75)$$

where $r_{\mathcal{BA}}(x,\epsilon)$ is the restriction of $(h(x,\epsilon) - E_{\mathcal{B}}(x,\epsilon)^{-1})$ to the range of $P_{\perp}^{\mathcal{AB}}(x)$.

The expression (3.73) is a standard leading order Born-Oppenheimer wave packet, and it suffices to show that it agrees with a solution of the Schrödinger equation up to errors on the order of ϵ^{ξ} . To prove this, we explicitly compute

$$\zeta(x,t,\epsilon) = i\epsilon^2 \frac{\partial\Psi}{\partial t} - H(\epsilon)\Psi.$$
(3.76)

For $-T \leq t \leq -\epsilon^{1-\xi}$, terms in $\zeta(x,t,\epsilon)$ that contain derivatives of F are easily seen to have norms that grow at worst like powers of 1/|t| times factors that are exponentially small in ϵ . Most of the remaining terms in $\zeta(x,t,\epsilon)$ are formally given on pages 105-108 of [16], with $(\frac{\partial}{\partial t} + \eta^{\mathcal{B}} \cdot \nabla_x)$ in place of $\eta^{\mathcal{B}} \cdot \nabla_x$ in many expressions. A few other terms arise from the $i\epsilon^2\frac{\partial}{\partial t}$ acting on the electronic eigenfunctions in (3.73). Our error term $\zeta(x,t,\epsilon)$ differs formally from that of Section 6 of [16] because our eigenfunctions have time dependence and the ones in Section 6 of [16] do not.

After noting that numerous terms in $\zeta(x, t, \epsilon)$ cancel with one another, we estimate the remaining terms individually. This process is very similar to that described on pages 108-110 of [16], except that there are a few extra terms that contain time derivatives of the electronic eigenfunctions. To estimate these terms we use arguments similar to those used in the proof of lemma 3.2. Our avoided crossing problem is slightly less singular than the crossing problem treated in [16]. So, it is not surprising that our error term satisfies the same estimate

$$\|\zeta(\cdot, t, \epsilon)\| \le C \epsilon^3 |t|^{-2} \tag{3.77}$$

as the corresponding error term in Section 6 of [16].

We obtain the desired estimate by using the estimate (3.77) and lemma 3.3 of [16]. \Box

3.4. The Inner Solution

In the inner time regime, we look for an approximation constructed by means of the classical quantities associated with the potential $\tilde{V}(x, \epsilon)$. Let a(t) and S(t) be the corresponding classical quantities satisfying the initial conditions

$$\begin{array}{c} a(0) = 0\\ \eta(0) = \eta^{0}\\ S(0) = 0. \end{array} \right\}$$
(3.78)

In the rescaled variables

$$\begin{cases} y = (x - a(t))/\epsilon, \\ s = t/\epsilon \end{cases}$$

$$(3.79)$$

the Schrödinger equation reads

$$i\epsilon \frac{\partial}{\partial s}\psi - i\epsilon\eta(\epsilon s)\nabla_y\psi = -\frac{\epsilon^2}{2}\Delta_y\psi + h(a(\epsilon s) + \epsilon y, \epsilon)\psi.$$
(3.80)

We seek an approximate solution of the form

$$F(\|y\|\epsilon^{\delta'}) \exp\left(i\frac{S(\epsilon s)}{\epsilon^2} + i\frac{\eta(\epsilon s)y}{\epsilon}\right)\chi(y,s,\epsilon),$$
(3.81)

with

$$\chi(y, s, \epsilon) = \{ f(y, s, \epsilon)\psi_1(a(\epsilon s) + \epsilon y, \epsilon) + g(y, s, \epsilon)\psi_2(a(\epsilon s) + \epsilon y, \epsilon) + \psi_{\perp}(a(\epsilon s) + \epsilon y, \epsilon) \},$$

$$(3.82)$$

where $\psi_{\perp}(x,\epsilon) \in (\mathbb{1} - P(x,\epsilon))\mathcal{H}$ and f, g are scalar functions. Anticipating exponential fall off of the solution, we insert (3.81) in (3.80) and neglect the

derivatives of F. Making use of the decomposition (1.11) and eliminating the overall "classical phase" and cut-off F, we get (dropping the arguments)

$$\begin{split} &i\epsilon\psi_{1}\frac{\partial}{\partial s}f + i\epsilon\psi_{2}\frac{\partial}{\partial s}g + i\epsilon\frac{\partial}{\partial s}\psi_{\perp} + (\widetilde{V}(a) + \epsilon y\nabla_{x}\widetilde{V}(a))(f\psi_{1} + g\psi_{2} + \psi_{\perp}) \\ &+ i\epsilon^{2}f\eta\nabla_{x}\psi_{1} + i\epsilon^{2}g\eta\nabla_{x}\psi_{2} + i\epsilon^{2}\eta\nabla_{x}\psi_{\perp} \\ &= -\epsilon^{4}\Delta_{x}\psi_{\perp}/2 - \epsilon^{4}f\Delta_{x}\psi_{1}/2 - \epsilon^{4}g\Delta_{x}\psi_{2}/2 - \epsilon^{3}\nabla_{y}f\nabla_{x}\psi_{1} \\ &- \epsilon^{3}\nabla_{y}g\nabla_{x}\psi_{2} - \epsilon^{2}\psi_{1}\Delta_{y}f - \epsilon^{2}\psi_{2}\Delta_{y}g + h^{\perp}\psi_{\perp} + h_{1}(f\psi_{1} + g\psi_{2}) \\ &+ \widetilde{V}(a + \epsilon y)(f\psi_{1} + g\psi_{2} + \psi_{\perp}). \end{split}$$
(3.83)

We assume the solutions have expansions of the form

$$f(y, s, \epsilon) = \sum_{j=0}^{\infty} \nu_j(\epsilon) f_j(y, s),$$

$$g(y, s, \epsilon) = \sum_{j=0}^{\infty} \nu_j(\epsilon) g_j(y, s),$$

$$\psi_{\perp}(x, t, \epsilon) = \sum_{j=0}^{\infty} \nu_j(\epsilon) \psi_{\perp j}(x, t),$$

$$(3.84)$$

with asymptotic scales $\nu_j(\epsilon)$ to be determined by matching. We insert these expansions in (3.83). Using the behaviors (1.20) and lemma 2.2 we verify that the lowest order terms yield

$$i\epsilon\nu_{0}(\epsilon)\frac{\partial}{\partial s}f_{0}(y,s)\psi_{1}(a(\epsilon s)+\epsilon y,\epsilon)+i\epsilon\nu_{0}(\epsilon)\frac{\partial}{\partial s}g_{0}(y,s)\psi_{2}(a(\epsilon s)+\epsilon y,\epsilon)$$

= $h_{\perp}(a(\epsilon s)+\epsilon y,\epsilon)\psi_{\perp}(a(\epsilon s)+\epsilon y,\epsilon s,\epsilon)$
+ $\epsilon\nu_{0}(\epsilon)h_{11}(y,s)(f_{0}(y,s)\psi_{1}(a(\epsilon s)+\epsilon y,\epsilon)+g_{0}(y,s)\psi_{2}(a(\epsilon s)+\epsilon y,\epsilon))$
(3.85)

on the support of F, where $h_{11}(y,s)$ is the operator on the span of $\{\psi_1(a(\epsilon s) + \epsilon y, \epsilon), \psi_2(a(\epsilon s) + \epsilon y, \epsilon)\}$ whose matrix in the basis $\{\psi_1(a(\epsilon s) + \epsilon y, \epsilon), \psi_2(a(\epsilon s) + \epsilon y, \epsilon)\}$ is given by

$$r \begin{pmatrix} \eta_1^0 s + y_1 & 1\\ 1 & -\eta_1^0 s - y_1 \end{pmatrix}.$$
 (3.86)

By matching the incoming outer solution, we obtain

$$\nu_j(\epsilon)h_{\perp}(a(\epsilon s) + \epsilon y, \epsilon)\psi_{\perp j}(a(\epsilon s) + \epsilon y, \epsilon s) = 0$$
(3.87)

if $\nu_j(\epsilon) << \epsilon \nu_0(\epsilon)$, $j = 0, 1, \dots, m-1$. The spectrum of $h_{\perp}(x, \epsilon)$ is bounded away from 0 in a neighborhood of (0, 0). This implies

$$\psi_{\perp j}(x,t) = 0, \quad j = 0, 1, \cdots, m-1.$$
 (3.88)

We split the remaining equation for the order $\epsilon \nu_0(\epsilon) = \nu_m(\epsilon)$ (without arguments)

$$i\epsilon\nu_{0}(\epsilon)\frac{\partial}{\partial s}f_{0}\psi_{1} + i\epsilon\nu_{0}(\epsilon)\frac{\partial}{\partial s}g_{0}\psi_{2}$$

= $\nu_{m}(\epsilon)h_{\perp}\psi_{\perp} + \epsilon\nu_{0}(\epsilon)h_{11}(f_{0}\psi_{1} + g_{0}\psi_{2})$ (3.89)

into

$$i\epsilon\nu_{0}(\epsilon)\frac{\partial}{\partial s}f_{0}\psi_{1} + i\epsilon\nu_{0}(\epsilon)\frac{\partial}{\partial s}g_{0}\psi_{2} = \epsilon\nu_{0}(\epsilon)h_{11}(f_{0}\psi_{1} + g_{0}\psi_{2}), \\ \nu_{m}(\epsilon)h_{\perp}\psi_{\perp m} = 0 \qquad (3.90)$$

by projection with $P(x, \epsilon)$ and $(1 - P(x, \epsilon))$. The second equation gives $\psi_{\perp m} = 0$. The first one is equivalent to

$$i\frac{\partial}{\partial s} \begin{pmatrix} f_0(y,s) \\ g_0(y,s) \end{pmatrix} = r \begin{pmatrix} \eta_1^0 s + y_1 & 1 \\ 1 & -(\eta_1^0 s + y_1) \end{pmatrix} \begin{pmatrix} f_0(y,s) \\ g_0(y,s) \end{pmatrix}.$$
(3.91)

The general solution of this equation can be found exactly in terms of parabolic cylinder functions [5].

$$\begin{pmatrix} f_{0}(y,s)\\ g_{0}(y,s) \end{pmatrix} = C_{1}(y) \begin{pmatrix} \frac{(1-i)}{2} \sqrt{\frac{r}{\eta_{1}^{0}}} D_{\frac{ir}{2\eta_{1}^{0}}-1} \left((-1+i) \sqrt{\frac{r}{\eta_{1}^{0}}} (\eta_{1}^{0}s+y_{1}) \right) \\ D_{\frac{ir}{2\eta_{1}^{0}}} \left((-1+i) \sqrt{\frac{r}{\eta_{1}^{0}}} (\eta_{1}^{0}s+y_{1}) \right) \\ + C_{2}(y) \begin{pmatrix} D_{-\frac{ir}{2\eta_{1}^{0}}} \left(-(1+i) \sqrt{\frac{r}{\eta_{1}^{0}}} (\eta_{1}^{0}s+y_{1}) \right) \\ -\frac{(1+i)}{2} \sqrt{\frac{r}{\eta_{1}^{0}}} D_{-\frac{ir}{2\eta_{1}^{0}}-1} \left(-(1+i) \sqrt{\frac{r}{\eta_{1}^{0}}} (\eta_{1}^{0}s+y_{1}) \right) \end{pmatrix}.$$

$$(3.92)$$

We define our inner approximation by taking

$$\psi_{I}(y,s) = F(||y||\epsilon^{\delta'}) \exp\left(i\frac{S(\epsilon s)}{\epsilon^{2}} + i\frac{\eta(\epsilon s)y}{\epsilon}\right)$$
$$\times [\nu_{0}(\epsilon)f_{0}(y,s)\psi_{1}(a(\epsilon s) + \epsilon y, \epsilon)$$
$$+ \nu_{0}(\epsilon)g_{0}(y,s)\psi_{2}(a(\epsilon s) + \epsilon y, \epsilon)], \qquad (3.93)$$

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with f_0 and g_0 as above. To get a vector along the \mathcal{B} level as t < 0, we need to impose

$$g_0(y,s) \rightarrow 0$$
, as $s \rightarrow -\infty$. (3.94)

Using the asymptotic formulae in [5] with $||y||/s = \mathcal{O}(\epsilon^{\xi - \delta'}) \rightarrow 0$, due to the cutoff, we have that

$$D_{\frac{ir}{2\eta_1^0}}\left((-1+i)\sqrt{\frac{r}{\eta_1^0}}(\eta_1^0s+y_1)\right) = e^{i\lambda(y,s)}e^{\frac{\pi r}{8\eta_1^0}}(1+\mathcal{O}(1/(s+||y||)^2))$$
(3.95)

as $s \rightarrow -\infty$, where $\lambda(y,s) \in \mathbb{R}$, so that

$$C_1(y) \equiv 0. \tag{3.96}$$

We further compute [5]

$$D_{-\frac{ir}{2\eta_1^0}} \left(-(1+i)\sqrt{\frac{r}{\eta_1^0}} (\eta_1^0 s + y_1) \right)$$

= $\exp\left\{ -\frac{ir}{2\eta_1^0} \left(\eta_1^{0^2} s^2 + 2y_1 \eta_1^0 s + y_1^2 + \ln|s| + \ln(2r\eta_1^0)/2 \right) \right\}$
 $\times e^{\frac{ir}{8\eta_1^0}} \left(1 + \mathcal{O}(||y||/s) + \mathcal{O}(1/s^2) \right), \text{ and}$
 $D_{-\frac{ir}{2\eta_1^0} - 1} \left(-(1+i)\sqrt{\frac{r}{\eta_1^0}} (\eta_1^0 s + y_1) \right) = \mathcal{O}(1/(s + ||y||)) \quad (3.97)$

as $s \to -\infty$ under the same conditions. We match $\psi_I(y,s)$ with $\psi_{IO}(x,t)$, the incoming outer solution of lemma 3.4, in the regime characterized by

$$t = \epsilon s, \ s = -\epsilon^{-\xi}, \ y = \mathcal{O}(\epsilon^{-\delta'}), \ x = \mathcal{O}(\epsilon^{\kappa})$$
 (3.98)

with

$$0 < \delta' < \xi < 1/3, \ 2/3 < \kappa < 1 - \xi.$$
(3.99)

Let us consider $\psi_{IO}(x,t)$ as $t, \epsilon \rightarrow 0$ and express the result in terms of the new variables (y,s). From $(a^{\mathcal{B}}(t) - a(t))/\epsilon = \mathcal{O}(t^2/\epsilon) = \mathcal{O}(\epsilon s)$, we deduce that

$$F(\|y^{\mathcal{B}}\|\epsilon^{\delta'}) = F(\|y\|\epsilon^{\delta'} + \mathcal{O}(\epsilon^{1+\delta'}s)) = F(\|y\|\epsilon^{\delta'}) + \mathcal{O}(\epsilon^{1+\delta'}s) \quad (3.100)$$

and from proposition 2.2, its corollary and lemma 3.1, we get, in the $L^2(\mathbb{R}^n)$ sense,

$$\varphi_{l}(A^{\mathcal{B}}(t), B^{\mathcal{B}}(t), \epsilon^{2}, a^{\mathcal{B}}(t), 0, x) = \epsilon^{-n/2} \varphi_{l}(A_{0}, B_{0}^{\mathcal{B}}(-), 1, 0, 0, y) + \mathcal{O}(\epsilon s) + \mathcal{O}(1/s^{2}).$$
(3.101)

We control the phases by means of the previous results on the asymptotic behaviors of the corresponding classical quantities. We use corollary 2.3 to see that

$$\exp\left\{iS^{\mathcal{B}}(\epsilon s)/\epsilon^{2}\right\}$$

$$=\exp\left\{iS(\epsilon s)/\epsilon^{2}+iS_{0}^{\mathcal{B}}(\epsilon,-)/\epsilon^{2}-irs-i\left(r\eta_{1}^{0}s^{2}+\frac{r}{\eta_{1}^{0}}\ln(\epsilon|s|)\right)\right\}$$

$$\times\left(1+\mathcal{O}(\epsilon s^{3})+\mathcal{O}(1/s^{2})+\mathcal{O}(\epsilon\ln(\epsilon|s|))\right),\qquad(3.102)$$

We use corollary 2.1 and corollary 2.2 to estimate

$$\exp\left\{i\eta^{\mathcal{B}}(\epsilon s)(x-a^{\mathcal{B}}(\epsilon s))/\epsilon^{2}\right\}$$

=
$$\exp\left\{i\eta(\epsilon s)y/\epsilon + i(\eta^{\mathcal{B}}(\epsilon s) - \eta(\epsilon s))y/\epsilon + i\eta^{\mathcal{B}}(\epsilon s)(a(\epsilon s) - a^{\mathcal{B}}(\epsilon s))/\epsilon^{2}\right\}$$

=
$$\exp\left\{i\eta(\epsilon s)y/\epsilon - iy_{1}(r/\eta_{1}^{0} + rs)\right\}$$

$$\times \exp\left\{irs + i\frac{r\eta_{1}^{0}}{2}s^{2} + i\frac{r}{2\eta_{1}^{0}}\left(\ln(\epsilon|s|) + \ln(2\eta_{1}^{0}) + 1/2 - \ln\epsilon\right)\right\}$$

$$\times \left(1 + \mathcal{O}(\epsilon s^{3}) + \mathcal{O}(1/s^{2}) + \mathcal{O}(||y||\epsilon s^{2}) + \mathcal{O}(||y||/s) + \mathcal{O}(\epsilon s \ln\epsilon)\right)$$

(3.103)

as $s \rightarrow -\infty$. Finally, by lemmas 3.2 and 3.3,

$$\Phi_{\mathcal{B}}^{-}(x,\epsilon s,\epsilon) = -\psi_1(x,\epsilon) + \mathcal{O}(\epsilon^{1/2}/s^{1/2}) + \mathcal{O}(s\epsilon^{\kappa})$$
(3.104)

so that we can do the matching provided $\nu_0(\epsilon) = 1$ and

$$C_{2}(y) = -\epsilon^{-n/2} \varphi_{l}(A_{0}, B_{0}^{\mathcal{B}}(-), 1, 0, 0, y) e^{-\frac{\pi r}{8\eta_{1}^{0}}} \exp\left(\frac{ir}{2\eta_{1}^{0}}(y_{1}^{2} - 2y_{1})\right) \\ \times \exp\left(i\frac{S_{0}^{\mathcal{B}}(\epsilon, -)}{\epsilon^{2}} + \frac{ir}{4\eta_{1}^{0}}(1 + 3\ln(2\eta_{1}^{0}) + \ln r - 4\ln\epsilon)\right).$$
(3.105)

Using the fact that $\varphi_l(A^{\mathcal{B}}(t), B^{\mathcal{B}}(t), \epsilon^2, a^{\mathcal{B}}(t), 0, x)$ is normalized, we see that the difference between these expressions in the regime defined by (3.99) is, in the $L^2(\mathbb{R}^n)$ sense,

$$\psi_{IO}(\epsilon y + a(\epsilon s), \epsilon s) = \psi_I(y, s) + \mathcal{O}(\epsilon^p)$$
(3.106)

for some positive p.

We need to check the validity of this inner approximation. We return to the initial variables and set

$$\Psi_I(x,t) = \psi_I((x-a(t))/\epsilon, t/\epsilon).$$
(3.107)

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If we define $\zeta(x,t)$ by

$$i\epsilon^2 \frac{\partial}{\partial t} \Psi_I(x,t) = -\frac{\epsilon^4}{2} \Delta \Psi_I(x,t) + h(x,\epsilon) \Psi_I(x,t) + \zeta(x,\epsilon), \quad (3.108)$$

then we have

PROPOSITION 3.1. – In the $L^2(\mathbb{R}^n)$ norm

$$\frac{1}{\epsilon^2} \int_{-\epsilon^{1-\xi}}^{\epsilon^{1-\xi}} \|\zeta(x,t)\| dt = \mathcal{O}(\epsilon^{1-3\xi}).$$

Thus, by virtue of lemma 3.3 in [16], we have,

LEMMA 3.5. – The function $\Psi_I(x,t)$ is an approximation of a solution $\psi(x,t)$ of the Schrödinger equation for $-\epsilon^{1-\xi} \leq t \leq \epsilon^{1-\xi}$, such that in the $L^2(\mathbb{R}^n)$ sense,

$$\psi(x,t) - \Psi_I(x,t) = \mathcal{O}(\epsilon^{1-3\xi}) \rightarrow 0,$$

as $\epsilon \rightarrow 0$.

Proof. – To compute $\zeta(x,t)$ we set

$$A(y,s) = F(||y||\epsilon^{\delta'})[f_0(y,s)\psi_1(a(\epsilon s) + \epsilon y, \epsilon) + g_0(y,s)\psi_2(a(\epsilon s) + \epsilon y, \epsilon)]$$
(3.109)

and define

$$\widetilde{\zeta}(y,s) = i\epsilon \frac{\partial}{\partial s} A(y,s) - H(y,s)A(y,s), \qquad (3.110)$$

where formally

$$\begin{split} H(y,s) &= -\frac{\epsilon^2}{2} \Delta_y + h_1(a(\epsilon s) + \epsilon y, \epsilon) \\ &+ \Big(\widetilde{V}(a(\epsilon s) + \epsilon y, \epsilon) - \widetilde{V}(a(\epsilon s), \epsilon) - \epsilon y \nabla_x \widetilde{V}(a(\epsilon s), \epsilon) \Big). \end{split}$$
(3.111)

It is readily checked that

$$\zeta(x,t) = \exp\left(i\frac{S(t)}{\epsilon^2} + i\frac{\eta(t)(x-a(t))}{\epsilon^2}\right)\widetilde{\zeta}((x-a(t))/\epsilon, t/\epsilon). \quad (3.112)$$

Then, making use of

$$\nabla_{y} F(\|y\|\epsilon^{\delta'}) = \epsilon^{\delta'} F'(\|y\|\epsilon^{\delta'})y/\|y\| \Delta_{y} F(\|y\|\epsilon^{\delta'}) = \epsilon^{2\delta'} F''(\|y\|\epsilon^{\delta'}),$$

$$(3.113)$$

the cancellations due to our choice of f_0 and g_0 , and the mean value theorem, we get (dropping the arguments)

$$\begin{split} \widetilde{\zeta}(y,s) &= \frac{1}{2} \epsilon^{2\delta'+2} F''(f_0 \psi_1 + g_0 \psi_2) \\ &+ \epsilon^{\delta'} F' \frac{y}{\|y\|} \left(\epsilon^2 \psi_1 \nabla_y f_0 + \epsilon^2 \psi_2 \nabla_y g_0 + \epsilon^3 f_0 \nabla_x \psi_1 + \epsilon^3 g_0 \nabla_x \psi_2 \right) \\ &+ F \left(i \epsilon^2 f_0 \eta \nabla_x \psi_1 + i \epsilon^2 g_0 \eta \nabla_x \psi_2 + \epsilon^3 \nabla_y f_0 \nabla_x \psi_1 + \epsilon^3 \nabla_y g_0 \nabla_x \psi_2 \right) \\ &+ \frac{\epsilon^2}{2} \psi_1 \Delta_y f_0 + \epsilon^2 / 2 \psi_2 \Delta_y g_0 + \epsilon^4 / 2 f_0 \Delta_y \psi_1 + \epsilon^4 / 2 g_0 \Delta_y \psi_2, \\ &- F \left(\psi_1, \psi_2 \right) \begin{pmatrix} \beta - \epsilon r(\eta_1^0 s + y_1) & \gamma + i \delta - r \epsilon \\ \gamma - i \delta - r \epsilon & -\beta + \epsilon r(\eta_1^0 s + y_1) \end{pmatrix} \begin{pmatrix} f_0 \\ g_0 \end{pmatrix} \\ &+ \frac{1}{2} \epsilon^2 y^2 F \widetilde{V}^{(2)} (a(\epsilon s) + \theta \epsilon y, \epsilon) (f_0 \psi_1 + g_0 \psi_2), \end{split}$$
(3.114)

where $\theta \in]0, 1[$. Then we use the facts that in the equation above we have • $||y|| = \mathcal{O}(\epsilon^{-\delta'}), s = \mathcal{O}(\epsilon^{-\xi}), \text{ with } \xi > \delta', t = \mathcal{O}(\epsilon^{1-\xi}), x = \epsilon y + a(t) = \mathcal{O}(\epsilon^{1-\xi}),$

• $F(||y||\epsilon^{\delta'}), F'(||y||\epsilon^{\delta'})$ and $F''(||y||\epsilon^{\delta'})$ are uniformly bounded in (y, ϵ) ,

• $\psi_j(x,\epsilon), \nabla_x \psi_j(x,\epsilon), \Delta_x \psi_j(x,\epsilon)$ and $\widetilde{V}^{(2)}(x,\epsilon)$ are uniformly bounded as (x,ϵ) goes to (0,0),

• $\eta(t)$ is uniformly bounded,

• $\beta - \epsilon r(\eta_1^0 s + y_1)$ and $\gamma - i\delta - r\epsilon$ are $\mathcal{O}(\epsilon^{2(1-\xi)})$, to get the estimate, for y in the support of F,

$$\begin{aligned} |\tilde{\zeta}(y,s)| &= \mathcal{O}\Big(\epsilon^{2(1-\xi)}(f_0(y,s) + g_0(y,s)) \\ &+ \epsilon^{2+\delta'}(\nabla_y f_0(y,s) + \nabla_y g_0(y,s)) + \epsilon^2(\triangle_y f_0(y,s) + \triangle_y g_0(y,s))\Big). \end{aligned}$$
(3.115)

We can write

$$\int_{-\epsilon^{1-\xi}}^{\epsilon^{1-\xi}} \|\zeta(\cdot,t)\|_{L^{2}(\mathbb{R}^{n})} dt$$

$$= \int_{-\epsilon^{1-\xi}}^{\epsilon^{1-\xi}} \left\{ \int_{\mathbb{R}^{n}} |\widetilde{\zeta}((x-a(t))/\epsilon,t/\epsilon)|^{2} dx \right\}^{1/2} dt$$

$$= \int_{-\epsilon^{-\xi}}^{\epsilon^{-\xi}} \epsilon^{n/2} \epsilon \|\widetilde{\zeta}(\cdot,s)\|_{L^{2}(\mathbb{R}^{n})} ds$$

$$\leq 2 \sup_{s \in [-\epsilon^{-\xi},\epsilon^{-\xi}]} \epsilon^{n/2} \|\widetilde{\zeta}(\cdot,s)\|_{L^{2}(\mathbb{R}^{n})} \epsilon^{1-\xi}.$$
(3.116)

Now f_0 , g_0 are given by a gaussian times a polynomial in y times a parabolic cylinder function. Asymptotically, these parabolic cylinder functions, their first and second derivatives are of order $\mathcal{O}((s + ||y||)^0)$, $\mathcal{O}((s + ||y||))$ and $\mathcal{O}((s + ||y||)^2)$, respectively, where $s = \mathcal{O}(\epsilon^{-\xi})$. Since the factor $\epsilon^{-n/2}$ in $C_2(y)$ takes care of the Jacobian, this yields the result. \Box

3.5 The Outgoing Outer Solution

Consider the classical quantities associated with the A level that satisfy the following initial conditions at t = 0:

$$a^{\mathcal{A}}(0) = 0$$
(3.117)

$$\eta^{\mathcal{A}}(0) = \eta^{0} - \epsilon \frac{2r}{\eta_{1}^{0}} \begin{pmatrix} 1\\0\\\vdots\\0 \end{pmatrix}$$

$$A^{\mathcal{A}}(0) = A_{0}$$

$$B^{\mathcal{A}}(0) = B_{0}$$

$$S^{\mathcal{A}}(0) = 0.$$
(3.118)

Note the difference in initial momentum with the classical quantities associated with the \mathcal{B} level, (3.69). As explained earlier, the corresponding loss in kinetic energy equals the gain in potential energy, to leading order. We look for an outgoing outer solution $\psi_{OO}(x,t)$, valid in the temporal region $\epsilon^{1-\xi} \leq t < T$ as a linear combination of standard Born-Oppenheimer wave packets. The first one is associated with the \mathcal{A} level and initial conditions (3.117), the other one is associated with the \mathcal{B} level and initial conditions (3.69). This function $\psi_{OO}(x,t)$ must match $\psi_I(y,s)$ in the temporal layer

$$t = \epsilon s, \ s = +\epsilon^{-\xi}, \ y = \mathcal{O}(\epsilon^{-\delta'}), \ x = \mathcal{O}(\epsilon^{\kappa})$$
 (3.119)

with (3.99). We proceed as above. We have [5]

$$\begin{split} &D_{-\frac{ir}{2\eta_1^0}} \left(-(1+i)\sqrt{\frac{r}{\eta_1^0}}(\eta_1^0 s+y_1) \right) \\ &= \exp\left\{ -\frac{ir}{2\eta_1^0} \left(\eta_1^{0^2} s^2 + 2y_1 \eta_1^0 s+y_1^2 + \ln|s| + \ln(2r\eta_1^0)/2 \right) \right\} \\ &\times \mathrm{e}^{-\frac{3r\pi}{8\eta_1^0}}(1+\mathcal{O}(||y||/s)) + \mathcal{O}(1/s), \\ &D_{-\frac{ir}{2\eta_1^0}-1} \left(-(1+i)\sqrt{\frac{r}{\eta_1^0}}(\eta_1^0 s+y_1) \right) \end{split}$$

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$$= \exp\left\{ + \frac{ir}{2\eta_1^0} \left(\eta_1^{0^2} s^2 + 2y_1 \eta_1^0 s + y_1^2 + \ln|s| + \ln(2r\eta_1^0)/2 \right) \right\}$$

$$\times \frac{\sqrt{2\pi}}{\Gamma\left(1 + \frac{ir}{2\eta_1^0}\right)} e^{-\frac{r\pi}{s\eta_1^0}} \left(1 + \mathcal{O}(1/s^2) \right) + \mathcal{O}(1/s), \qquad (3.120)$$

as $s \to +\infty$ and $||y||/s = \mathcal{O}(\epsilon^{\xi - \delta'}) \to 0$. Equation (3.100) is also satisfied if we replace the index \mathcal{B} by \mathcal{A} and since $B_0^{\mathcal{B}}(-) \equiv B_0^{\mathcal{A}}(+)$ (see corollary 2.4), we have

$$\varphi_l(A_0, B_0^{\mathcal{B}}(-), 1, 0, 0, y) \equiv \varphi_l(A_0, B_0^{\mathcal{A}}(+), 1, 0, 0, y).$$
(3.121)

On the other hand, using the notation (3.24), we can write

$$\varphi_{l}(A_{0}, B_{0}^{\mathcal{B}}(-), 1, 0, 0, y) = h(A_{0}, |A_{0}|^{-1}y) \exp\left\{-\frac{\langle y|B_{0}A_{0}^{-1}y\rangle}{2} - i\frac{\langle y|rPy\rangle}{2\eta_{1}^{0}}\right\} = \varphi_{l}(A_{0}, B_{0}^{\mathcal{B}}(+), 1, 0, 0, y)e^{-iy_{1}^{2}r/\eta_{1}^{0}}.$$
(3.122)

The equivalent of equation (3.101) for $s \rightarrow +\infty$ and indices A and B is satisfied with the same error terms. For the phases we have by corollary 2.3,

$$\exp\left\{iS^{\mathcal{B}}(\epsilon s)/\epsilon^{2}\right\}$$

$$=\exp\left\{iS(\epsilon s)/\epsilon^{2}+iS_{0}^{\mathcal{B}}(\epsilon,+)/\epsilon^{2}-irs+i\left(r\eta_{1}^{0}s^{2}+\frac{r}{\eta_{1}^{0}}\ln(\epsilon s)\right)\right\}$$

$$\times\left(1+\mathcal{O}(\epsilon s^{3})+\mathcal{O}(1/s^{2})+\mathcal{O}(\epsilon\ln(\epsilon s))\right),$$
(3.123)

and

$$\exp\left\{iS^{\mathcal{A}}(\epsilon s)/\epsilon^{2}\right\}$$

$$=\exp\left\{iS(\epsilon s)/\epsilon^{2}+iS_{0}^{\mathcal{A}}(\epsilon,+)/\epsilon^{2}-irs-i\left(r\eta_{1}^{0}s^{2}+\frac{r}{\eta_{1}^{0}}\ln(\epsilon|s|)\right)\right\}$$

$$\times\left(1+\mathcal{O}(\epsilon s^{3})+\mathcal{O}(1/s^{2})+\mathcal{O}(\epsilon\ln(\epsilon|s|))\right).$$
(3.124)

By corollary 2.1 and corollary 2.2,

$$\exp\left\{i\eta^{\mathcal{B}}(\epsilon s)(x-a^{\mathcal{B}}(\epsilon s))/\epsilon^{2}\right\}$$

$$=\exp\left\{i\eta(\epsilon s)y/\epsilon - iy_{1}(r/\eta_{1}^{0}-rs)\right\}$$

$$\times \exp\left\{irs - i\frac{r\eta_{1}^{0}}{2}s^{2} - i\frac{r}{2\eta_{1}^{0}}\left(\ln(\epsilon s) + \ln(2\eta_{1}^{0}) + 1/2 - \ln\epsilon\right)\right\}$$

$$\times \left(1 + \mathcal{O}(\epsilon s^{3}) + \mathcal{O}(1/s^{2}) + \mathcal{O}(||y||\epsilon s^{2}) + \mathcal{O}(||y||/s) + \mathcal{O}(\epsilon s \ln\epsilon)\right)$$

$$(3.125)$$

$$\exp\left\{i\eta^{\mathcal{A}}(\epsilon s)(x-a^{\mathcal{A}}(\epsilon s))/\epsilon^{2}\right\}$$

$$=\exp\left\{i\eta(\epsilon s)y/\epsilon - iy_{1}(r/\eta_{1}^{0}+rs)\right\}$$

$$\times \exp\left\{irs + i\frac{r\eta_{1}^{0}}{2}s^{2} + i\frac{r}{2\eta_{1}^{0}}\left(\ln(\epsilon s) + \ln(2\eta_{1}^{0}) + 1/2 - \ln\epsilon\right)\right\}$$

$$\times \left(1 + \mathcal{O}(\epsilon s^{3}) + \mathcal{O}(1/s^{2}) + \mathcal{O}(||y||/s) + \mathcal{O}(\epsilon s \ln\epsilon)\right)$$
(3.126)

as $s \rightarrow +\infty$. By lemmas 3.2 and 3.3

$$\Phi_{\mathcal{A}}^{+}(x,\epsilon s,\epsilon) = \psi_{1}(x,\epsilon) + \mathcal{O}(\epsilon^{1/2}/s^{1/2}) + \mathcal{O}(s\epsilon^{\kappa}),$$

$$\Phi_{\mathcal{B}}^{+}(x,\epsilon s,\epsilon) = \psi_{2}(x,\epsilon) + \mathcal{O}(\epsilon^{1/2}/s^{1/2}) + \mathcal{O}(s\epsilon^{\kappa}).$$
(3.127)

By matching ψ_I and ψ_{OO} in the domain defined by (3.119) we get up to terms of order $\mathcal{O}(\epsilon^p)$ in the $L^2(\mathbb{R}^n)$ sense

$$\begin{split} \psi_{OO}(x,t) &= - \,\mathrm{e}^{-\pi r/2\eta_{1}^{0}} F(\|x-a^{\mathcal{A}}(t)\|/\epsilon^{1-\delta'}) \mathrm{e}^{\left(i\frac{S^{\mathcal{A}}(t)}{\epsilon^{2}}+i\frac{\eta^{\mathcal{A}}(t)(x-a^{\mathcal{A}}(t))}{\epsilon^{2}}\right)} \\ &\qquad \times \varphi_{l}(A^{\mathcal{A}}(t), B^{\mathcal{A}}(t), \epsilon^{2}, a^{\mathcal{A}}(t), 0, x) \Phi^{+}_{\mathcal{A}}(x, t, \epsilon) \\ &+ \,\mathrm{e}^{-\pi r/4\eta_{1}^{0}} \sqrt{\frac{\pi r}{\eta_{1}^{0}}} \frac{e^{i\lambda(\epsilon)}}{\Gamma\left(1+\frac{ir}{2\eta_{1}^{0}}\right)} \\ &\qquad \times F(\|x-a^{\mathcal{B}}(t)\|/\epsilon^{1-\delta'}) \mathrm{e}^{\left(i\frac{S^{\mathcal{B}}(t)}{\epsilon^{2}}+i\frac{\eta^{\mathcal{B}}(t)(x-a^{\mathcal{B}}(t))}{\epsilon^{2}}\right)} \\ &\qquad \times \varphi_{l}(A^{\mathcal{B}}(t), B^{\mathcal{B}}(t), \epsilon^{2}, a^{\mathcal{B}}(t), 0, x) \Phi^{+}_{\mathcal{B}}(x, t, \epsilon), \quad (3.128) \end{split}$$

where

$$\lambda(\epsilon) = \pi/4 + S_0^{\mathcal{A}}(-,\epsilon)/\epsilon^2 + \frac{r}{2\eta_1^0}(1+3\ln(2\eta_1^0) + \ln r - 4\ln\epsilon). \quad (3.129)$$

We show as above that this function is an approximate solution of the Schrödinger equation.

LEMMA 3.6. – For any $\epsilon^{1-\xi} \leq t \leq T$, the function $\psi_{OO}(x,t)$ is an approximation of a solution of the Schrödinger equation $\psi(x,t)$, such that

$$\psi(x,t) = \psi_{OO}(x,t) + R(x,t,\epsilon)$$
 (3.130)

with

$$R(x.t.\epsilon) = \mathcal{O}(\epsilon^{\xi}) \tag{3.331}$$

in the $L^2(\mathbb{R}^n)$ sense, as $\epsilon \rightarrow 0$.

Proof. – This lemma is proved by the same techniques used to prove lemma 3.4. \Box

Remark. – In lemmas 3.4, 3.5 and 3.6, the cut-off functions F appearing in the approximations can be dropped without altering the estimates because of the exponential decay of the gaussians they multiply.

Gathering the results of this last section and using the fact that the evolution operator is unitary, we obtain our main theorem:

THEOREM 3.1. – Let $h(x, \epsilon)$ be a hamiltonian such that $h_{\parallel}(x, \epsilon)$ is characterized by (1.11) and (1.20), and let $\psi(x, t)$ be a solution of the corresponding Schrödinger equation (1.15) such that

$$\psi(x, -T) = e^{\left(i\frac{S^{\mathcal{B}}(-T)}{\epsilon^{2}} + i\frac{\eta^{\mathcal{B}}(-T)(x-a^{\mathcal{B}}(-T))}{\epsilon^{2}}\right)} \times \varphi_{l}(A^{\mathcal{B}}(-T), B^{\mathcal{B}}(-T), \epsilon^{2}, a^{\mathcal{B}}(-T), 0, x)\Phi_{\mathcal{B}}^{-}(x, -T, \epsilon) + \mathcal{O}(\epsilon^{q})$$

for some positive q, in the $L^2(\mathbb{R}^n)$ sense. Then, for any $0 < \xi < 1/3$, there exists a positive p such that in the limit $\epsilon \rightarrow 0$ we have for $-T \leq t \leq -\epsilon^{1-\xi}$:

$$\psi(x,t) = e^{\left(i\frac{\mathcal{S}^{\mathcal{B}}(t)}{\epsilon^{2}} + i\frac{\eta^{\mathcal{B}}(t)(x-a^{\mathcal{B}}(t))}{\epsilon^{2}}\right)} \times \varphi_{l}(A^{\mathcal{B}}(t), B^{\mathcal{B}}(t), \epsilon^{2}, a^{\mathcal{B}}(t), 0, x)\Phi_{\mathcal{B}}^{-}(x, t, \epsilon) + \mathcal{O}(\epsilon^{p}),$$

and for $-\epsilon^{1-\xi} \leq t \leq \epsilon^{1-\xi}$:

$$\psi(x,t) = e^{\left(i\frac{S(t)}{\epsilon^2} + i\frac{\eta(t)(x-a(t))}{\epsilon^2}\right)} (f_0(y,s)\psi_1(x,\epsilon) + g_0(y,s)\psi_2(x,\epsilon)),$$

where f_0 and g_0 are defined in (3.92), (3.96) and (3.105) with $y = (x - a(t))/\epsilon$. $s = t/\epsilon$, and finally for $\epsilon^{1-\xi} \leq t \leq T$:

$$\begin{split} \psi(x,t) &= -e^{-\pi r/2\eta_1^0} e^{\left(i\frac{\mathcal{S}^{\mathcal{A}(t)}}{\epsilon^2} + i\frac{\eta^{\mathcal{A}(t)}(x-a^{\mathcal{A}(t)})}{\epsilon^2}\right)} \\ &\times \varphi_l(A^{\mathcal{A}}(t), B^{\mathcal{A}}(t), \epsilon^2, a^{\mathcal{A}}(t), 0, x) \Phi^+_{\mathcal{A}}(x, t, \epsilon) \\ &+ e^{-\pi r/4\eta_1^0} \sqrt{\frac{\pi r}{\eta_1^0}} \frac{e^{i\lambda(\epsilon)}}{\Gamma\left(1 + \frac{ir}{2\eta_1^0}\right)} \\ &\times e^{\left(i\frac{\mathcal{S}^{\mathcal{B}(t)}}{\epsilon^2} + i\frac{\eta^{\mathcal{B}(t)}(x-a^{\mathcal{B}}(t))}{\epsilon^2}\right)} \\ &\times \varphi_l(A^{\mathcal{B}}(t), B^{\mathcal{B}}(t), \epsilon^2, a^{\mathcal{B}}(t), 0, x) \Phi^+_{\mathcal{B}}(x, t, \epsilon) + \mathcal{O}(\epsilon^p). \end{split}$$

where $\lambda(\epsilon)$ is given in (3.129).

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As a direct corollary, we compute that the transition probability $\mathcal{P}_{\mathcal{B}\to\mathcal{A}}$ from the level \mathcal{B} at t = -T to the level \mathcal{A} at t = +T is given in the limit $\epsilon \to 0$ by

$$\mathcal{P}_{\mathcal{B}\to\mathcal{A}} = \mathrm{e}^{-\pi r/\eta_1^0} + \mathcal{O}(\epsilon^p). \tag{3.132}$$

As explained above, $\mathcal{P}_{\mathcal{B}\to\mathcal{A}}$ coincides with the answer provided by the Landau-Zener formula where the gap is given by ϵ . Of course, the probability to remain on the \mathcal{B} level satisfies

$$\mathcal{P}_{\mathcal{B}\to\mathcal{B}} = e^{-\pi r/2\eta_1^0} \frac{\pi r}{\eta_1^0 \left| \Gamma\left(1 + \frac{ir}{2\eta_1^0}\right) \right|^2} + \mathcal{O}(\epsilon^p)$$

= 1 - e^{-\pi r/\eta_1^0} + \mathcal{O}(\epsilon^p). (3.133)

Remark. – The above formulae coincide with (1.2) and (1.3) when one takes into account the change of variables performed at the end of the first section. Indeed, it is not difficult to verify that in terms of the initial variables $(x,t) \in \mathbb{R}^{n+1}$ the classical position f(t) and momentum $\mu(t)$ are related to the classical position a(t') and momentum $\eta(t')$ associated with the new variables $(x',t') \in \mathbb{R}^{n+1}$ (1.16) by

$$f(t) = a(t')/b_1 - \begin{pmatrix} b_2 \epsilon/b_1 \\ 0 \\ \vdots \\ 0 \end{pmatrix}$$

$$\mu(t) = \eta(t')b_1/c_2^2, \qquad (3.134)$$

which yields the result.

4. Type 2 Avoided Crossings

In this section we make a precise statement of our result for Type 2 Avoided Crossings and describe how to alter the arguments of Section 3 to produce a proof. The theorem is proved by the same arguments as Theorem 3.1 except for the adiabatic motion of the degenerate electronic states.

We choose electronic eigenvectors by mimicking Section 3.2. We let $\eta^{c}(t)$ be the solution of the classical equations of motion (2.2) and (2.3). The normalized eigenvectors, $\Phi_{C,1}^{\pm}(x,t,\epsilon)$ and $\Phi_{C,2}^{\pm}(x,t,\epsilon)$, we wish to construct are the orthonormal solutions of

$$\langle \Phi^{\pm}_{\mathcal{C},j}(x,t,\epsilon) | \big(\partial/\partial t + \eta^{\mathcal{C}}(t) \nabla \big) \Phi^{\pm}_{\mathcal{C},k}(x,t,\epsilon) \rangle \equiv 0$$
(4.1)

for C = A, B, j = 1, 2, k = 1, 2, and $t \ge 0$. Since the eigenvalues $E_A(x, \epsilon)$ and $E_B(x, \epsilon)$ are multiplicity 2 for any time t, t small enough, such vectors exist, are unique up to an overall time-dependent unitary transformation, and are eigenvectors of $h_1(x, \epsilon)$ associated with $E_C(x, \epsilon)$ for any time, see [16]. Let us make the construction of these eigenvectors more specific and give their small t and small ϵ asymptotics.

As in Section 3.2, we define $\varphi(x, \epsilon)$ and $\theta(x, \epsilon)$ by equations (3.33)-(3.35), and construct the following static eigenvectors:

$$\Phi_{\mathcal{A},1}^{-}(x,\epsilon) = e^{i\varphi(x,\epsilon)}\cos(\theta(x,\epsilon)/2)\psi_{1}(x,\epsilon) + \sin(\theta(x,\epsilon)/2)\psi_{2}(x,\epsilon), \qquad (4.2)$$

$$\Phi_{\mathcal{B},1}^{-}(x,\epsilon) = e^{-i\varphi(x,\epsilon)}\cos(\theta(x,\epsilon)/2)\psi_{2}(x,\epsilon) -\sin(\theta(x,\epsilon)/2)\psi_{1}(x,\epsilon),$$
(4.3)

$$\Phi_{\mathcal{A},2}^{-}(x,\epsilon) = e^{i\varphi(x,\epsilon)}\cos(\theta(x,\epsilon)/2)\psi_{3}(x,\epsilon)$$
(4.4)

$$+\sin(\theta(x,\epsilon)/2)\psi_4(x,\epsilon), \qquad (4.4)$$

$$\Phi_{\mathcal{B},2}^{-}(x,\epsilon) = e^{-i\varphi(x,\epsilon)}\cos(\theta(x,\epsilon)/2)\psi_4(x,\epsilon) -\sin(\theta(x,\epsilon)/2)\psi_3(x,\epsilon),$$
(4.5)

associated with $E_{\mathcal{C}}(x,\epsilon)$, $\mathcal{C} = \mathcal{A}, \mathcal{B}$, for $\pi/2 < \theta(x,\epsilon) \leq \pi$, and

$$\Phi_{\mathcal{A},1}^{+}(x,\epsilon) = \cos(\theta(x,\epsilon)/2)\psi_{1}(x,\epsilon) + e^{-i\varphi(x,\epsilon)}\sin(\theta(x,\epsilon)/2)\psi_{2}(x,\epsilon),$$
(4.6)

$$\Phi_{\mathcal{B},1}^+(x,\epsilon) = \cos(\theta(x,\epsilon)/2)\psi_2(x,\epsilon)$$

$$-e^{i\varphi(x,\epsilon)}\sin(\theta(x,\epsilon)/2)\psi_1(x,\epsilon), \qquad (4.7)$$

$$\Phi^+_{\mathcal{A},2}(x,\epsilon) = \cos(\theta(x,\epsilon)/2)\psi_3(x,\epsilon)$$

$$+ e^{-i\varphi(x,\epsilon)}\sin(\theta(x,\epsilon)/2)\psi_4(x,\epsilon), \qquad (4.8)$$

$$\Phi^{+}_{\mathcal{B},2}(x,\epsilon) = \cos(\theta(x,\epsilon)/2)\psi_{4}(x,\epsilon) - e^{i\varphi(x,\epsilon)}\sin(\theta(x,\epsilon)/2)\psi_{3}(x,\epsilon),$$
(4.9)

for $0 \le \theta(x, \epsilon) < \pi/2$. We then construct solutions of (4.1) by using

$$\begin{pmatrix} \Phi_{\mathcal{C},1}^{\pm}(x,t,\epsilon) \\ \Phi_{\mathcal{C},2}^{\pm}(x,t,\epsilon) \end{pmatrix} = U_{\mathcal{C}}^{\pm}(x,t,\epsilon) \begin{pmatrix} \Phi_{\mathcal{C},1}(x,\epsilon) \\ \Phi_{\mathcal{C},2}(x,\epsilon) \end{pmatrix}, \quad \begin{cases} t > 0 \\ t < 0 \end{cases}, \quad (4.10)$$

where the unitary matrices $U_{\mathcal{C}}^{\pm}(x,t,\epsilon)$ play the roles of the phase factors from Section 3.2. The matrices $U_{\mathcal{C}}^{\pm}(x,t,\epsilon)$ satisfy differential equations that lead to time ordered exponentials instead of simple exponentials, so

there is not a simple analog of equation (3.41). However, the analog of Lemma 3.2 is true.

LEMMA 4.1. – Let $\eta^{\mathcal{C}}(t)$ be the momentum solution of the classical equations of motion (2.2) and (2.3). There exist eigenvectors $\Phi_{\mathcal{C},j}^{\pm}(x,t,\epsilon)$, $\mathcal{C} = \mathcal{A}.\mathcal{B}, j = 1, 2$, such that

$$\langle \Phi^{\pm}_{\mathcal{C},j}(x,t,\epsilon) | (\partial/\partial t + \eta^{\mathcal{C}}(t)\nabla) \Phi^{\pm}_{\mathcal{C},k}(x,t,\epsilon) \rangle \equiv 0.$$

If $||x|| = \mathcal{O}(\epsilon^{\kappa})$, with $2/3 < \kappa < 1$,

$$\left|\Phi_{\mathcal{C},j}^{\pm}(x,t,\epsilon) - \Phi_{\mathcal{C},j}^{\pm}(x,\epsilon)\right| = \mathcal{O}(t/\epsilon^{1-\kappa}).$$

Proof. – We give a proof for $\Phi_{\mathcal{A},j}^+(x,t,\epsilon)$ only. The other cases are similar. In order to simplify the notation, we drop the indices \mathcal{A} and +. Thus we need to determine $U(x,t,\epsilon)$ such that

$$\Phi_j(x,t,\epsilon) = U(x,t,\epsilon) \Phi_j(x,\epsilon)$$
(4.11)

for j = 1, 2, where

$$\Phi_1(x,\epsilon) = \cos(\theta(x,\epsilon)/2)\psi_1(x,\epsilon) + e^{-i\varphi(x,\epsilon)}\sin(\theta(x,\epsilon)/2)\psi_2(x,\epsilon)$$
(4.12)

and

$$\Phi_2(x,\epsilon) = \cos(\theta(x,\epsilon)/2)\psi_3(x,\epsilon) + e^{-i\varphi(x,\epsilon)}\sin(\theta(x,\epsilon)/2)\psi_4(x,\epsilon)$$
(4.13)

Let $a(t) \equiv a^{\mathcal{A}}(t)$, and $\eta(t) \equiv \eta^{\mathcal{A}}(t)$ for t > 0, and introduce

$$\omega \equiv x - a(t), \tag{4.14}$$

$$\widetilde{U}(\omega,t,\epsilon)\equiv U(\omega+a(t),t,\epsilon),\qquad {\rm and}\qquad (4.15)$$

$$\widetilde{\Phi_j}(\omega, t, \epsilon) \equiv \Phi_j(\omega + a(t), \epsilon).$$
(4.16)

In terms of these new variables, the equation for the entries of \widetilde{U} has the form

$$i\frac{\partial}{\partial t}\widetilde{U}(\omega,t,\epsilon) = M(\omega,t,\epsilon)\widetilde{U}(\omega,t,\epsilon), \qquad (4.17)$$

where the precise form for $M(\omega, t, \epsilon)$ can be obtained from page 36 of [16], with

$$\frac{\partial}{\partial t}\widetilde{\Phi}_j(\omega, t, \epsilon) = \eta(t)\nabla\Phi_j(\omega + a(t), \epsilon).$$
(4.18)

In analogy with the computations of Section 3.2, we compute

$$\nabla \Phi_1 = \left(-\sin(\theta/2)\psi_1 + e^{-i\varphi}\cos(\theta/2)\psi_2\right)\nabla\theta/2 + \cos(\theta/2)\nabla\psi_1 + e^{-i\varphi}\sin(\theta/2)\nabla\psi_2 - i\nabla\varphi e^{-i\varphi}\sin(\theta/2)\psi_2, \qquad (4.19)$$

and

$$\nabla \Phi_2 = \left(-\sin(\theta/2)\psi_3 + e^{-i\varphi}\cos(\theta/2)\psi_4 \right) \nabla \theta/2 + \cos(\theta/2)\nabla\psi_3 + e^{-i\varphi}\sin(\theta/2)\nabla\psi_4 - i\nabla\varphi e^{-i\varphi}\sin(\theta/2)\psi_4.$$
(4.20)

From this we can compute the matrix elements of $M(\omega, t, \epsilon)$, which are the analogs of (3.50). The terms of the form

$$\left\langle \widetilde{\Phi_{1}} \middle| \frac{\partial}{\partial t} \widetilde{\Phi_{1}} \right\rangle = -i\eta \nabla \varphi \sin^{2}(\theta/2) + \cos^{2}(\theta/2) \langle \psi_{1} \middle| \eta \nabla \psi_{1} \rangle + \sin^{2}(\theta/2) \langle \psi_{2} \middle| \eta \nabla \psi_{2} \rangle + \sin(\theta/2) \cos(\theta/2) \left(e^{-i\varphi} \langle \psi_{1} \middle| \eta \nabla \psi_{2} \rangle + e^{i\varphi} \langle \psi_{2} \middle| \eta \nabla \psi_{1} \rangle \right)$$

$$(4.21)$$

and

$$\left\langle \widetilde{\Phi_{2}} \middle| \frac{\partial}{\partial t} \widetilde{\Phi_{2}} \right\rangle = -i\eta \nabla \varphi \sin^{2}(\theta/2) + \cos^{2}(\theta/2) \langle \psi_{3} | \eta \nabla \psi_{1} \rangle + \sin^{2}(\theta/2) \langle \psi_{4} | \eta \nabla \psi_{4} \rangle + \sin(\theta/2) \cos(\theta/2) \left(e^{-i\varphi} \langle \psi_{3} | \eta \nabla \psi_{4} \rangle + e^{i\varphi} \langle \psi_{4} | \eta \nabla \psi_{3} \rangle \right)$$

$$(4.22)$$

both satisfy the estimates from the proof of Lemma 3.2. Because $\{\psi_1, \psi_2, \psi_3, \psi_4\}$ are mutually orthogonal, there is considerable simplification of the terms

$$\left\langle \widetilde{\Phi_{1}} | \frac{\partial}{\partial t} \widetilde{\Phi_{2}} \right\rangle = \left\langle \cos(\theta/2) \psi_{1} + e^{-i\varphi} \sin(\theta/2) \psi_{2} | \cos(\theta/2) \nabla \psi_{3} + e^{-i\varphi} \sin(\theta/2) \nabla \psi_{4} \right\rangle$$

$$(4.23)$$

and

$$\left\langle \widetilde{\Phi_{2}} | \frac{\partial}{\partial t} \widetilde{\Phi_{1}} \right\rangle = \left\langle \cos(\theta/2) \psi_{3} + e^{-i\varphi} \sin(\theta/2) \psi_{4} | \cos(\theta/2) \nabla \psi_{1} + e^{-i\varphi} \sin(\theta/2) \nabla \psi_{2} \right\rangle.$$

$$(4.24)$$

From these formulas, we easily see that these terms have no singular behavior for small t and fixed $\epsilon > 0$. From the convergent Dyson series for $\tilde{U}(\omega, t, \epsilon)$ from equation (4.17) with initial condition $\tilde{U}(\omega, 0, \epsilon) = 1$, we then conclude that

$$\widetilde{U}(\omega, t, \epsilon) = \widetilde{U}(\omega, 0, \epsilon) + \mathcal{O}(t/\epsilon^{1-\kappa}).$$
(4.25)

This implies the lemma. \Box

The following analog of Lemma 3.3 is proved by the same techniques we used to prove Lemma 3.3.

LEMMA 4.2. – For x in the support of $F(||x - a^{\mathcal{C}}(t)||/\epsilon^{1-\delta'})$, $\mathcal{C} = \mathcal{A}, \mathcal{B}$, we have for t < 0,

$$\begin{split} \Phi_{\mathcal{A},1}^{-}(x,\epsilon) &= \psi_2(x,\epsilon) + \mathcal{O}(\epsilon/t^{1/2}), \\ \Phi_{\mathcal{B},1}^{-}(x,\epsilon) &= -\psi_1(x,\epsilon) + \mathcal{O}(\epsilon/t^{1/2}), \\ \Phi_{\mathcal{A},2}^{-}(x,\epsilon) &= \psi_4(x,\epsilon) + \mathcal{O}(\epsilon/t^{1/2}), \\ \Phi_{\mathcal{B},2}^{-}(x,\epsilon) &= -\psi_3(x,\epsilon) + \mathcal{O}(\epsilon/t^{1/2}), \end{split}$$

and for t > 0,

$$\begin{split} \Phi^+_{\mathcal{A},1}(x,\epsilon) &= \psi_1(x,\epsilon) + \mathcal{O}(\epsilon/t^{1/2}), \\ \Phi^+_{\mathcal{B},1}(x,\epsilon) &= \psi_2(x,\epsilon) + \mathcal{O}(\epsilon/t^{1/2}), \\ \Phi^+_{\mathcal{A},2}(x,\epsilon) &= \psi_3(x,\epsilon) + \mathcal{O}(\epsilon/t^{1/2}), \\ \Phi^+_{\mathcal{B},2}(x,\epsilon) &= \psi_4(x,\epsilon) + \mathcal{O}(\epsilon/t^{1/2}). \end{split}$$

Our main result of this section is the following:

THEOREM 4.1. – Let $h(x, \epsilon)$ be a hamiltonian such that $h_{\parallel}(x, \epsilon)$ is characterized by (1.14) and (1.20). Let $\psi(x, t)$ be a solution of the corresponding Schrödinger equation (1.15) such that

$$\psi(x, -T) = e^{\left(i\frac{\mathcal{B}^{\mathcal{B}}(-T)}{\epsilon^{2}} + i\frac{\eta^{\mathcal{B}}(-T)(x-a^{\mathcal{B}}(-T))}{\epsilon^{2}}\right)} \times \varphi_{l}(A^{\mathcal{B}}(-T), B^{\mathcal{B}}(-T), \epsilon^{2}, a^{\mathcal{B}}(-T), 0, x)\Phi_{\mathcal{B},1}^{-}(x, -T, \epsilon) + \mathcal{O}(\epsilon^{q}),$$

for some positive q, in the $L^2(\mathbb{R}^n)$ sense. Then, for any $0 < \xi < 1/3$, there exists a positive p such that in the limit $\epsilon \rightarrow 0$ we have for $-T \leq t \leq -\epsilon^{1-\xi}$:

$$\psi(x,t) = e^{\left(i\frac{\mathcal{S}^{\mathcal{B}}(t)}{\epsilon^{2}} + i\frac{\eta^{\mathcal{B}}(t)(x-a^{\mathcal{B}}(t))}{\epsilon^{2}}\right)} \times \varphi_{l}(A^{\mathcal{B}}(t), B^{\mathcal{B}}(t), \epsilon^{2}, a^{\mathcal{B}}(t), 0, x)\Phi_{\mathcal{B},1}^{-}(x,t,\epsilon) + \mathcal{O}(\epsilon^{p}),$$

and for $-\epsilon^{1-\xi} \leq t \leq \epsilon^{1-\xi}$:

$$\psi(x,t) = e^{\left(i\frac{S(t)}{\epsilon^2} + i\frac{\eta(t)(x-a(t))}{\epsilon^2}\right)} (f_0(y,s)\psi_1(x,\epsilon) + g_0(y,s)\psi_2(x,\epsilon)),$$

where f_0 and g_0 are defined in (3.92), (3.96) and (3.105) with $y = (x - a(t))/\epsilon$, $s = t/\epsilon$, and finally for $\epsilon^{1-\xi} \leq t \leq T$:

$$\begin{split} \psi(x,t) &= -e^{-\pi r/2\eta_1^0} e^{\left(i\frac{S^{\mathcal{A}}(t)}{\epsilon^2} + i\frac{\eta^{\mathcal{A}}(t)(x-a^{\mathcal{A}}(t))}{\epsilon^2}\right)} \\ &\times \varphi_l(A^{\mathcal{A}}(t), B^{\mathcal{A}}(t), \epsilon^2, a^{\mathcal{A}}(t), 0, x) \Phi^+_{\mathcal{A},1}(x, t, \epsilon) \\ &+ e^{-\pi r/4\eta_1^0} \sqrt{\frac{\pi r}{\eta_1^0}} \frac{e^{i\lambda(\epsilon)}}{\Gamma\left(1 + \frac{ir}{2\eta_1^0}\right)} \\ &\times e^{\left(i\frac{S^{\mathcal{B}}(t)}{\epsilon^2} + i\frac{\eta^{\mathcal{B}}(t)(x-a^{\mathcal{B}}(t))}{\epsilon^2}\right)} \\ &\times \varphi_l(A^{\mathcal{B}}(t), B^{\mathcal{B}}(t), \epsilon^2, a^{\mathcal{B}}(t), 0, x) \Phi^+_{\mathcal{B},1}(x, t, \epsilon) + \mathcal{O}(\epsilon^p), \end{split}$$

where $\lambda(\epsilon)$ is given in (3.129). The corresponding statement with $\Phi_{B,2}^-$ in place of $\Phi_{B,1}^-$ is true with ψ_1 and ψ_2 replaced by ψ_3 and ψ_4 , respectively.

The proof of this theorem is essentially the same as that for Theorem 3.1. There is no change in the semiclassical mechanics, and the lemmas of this section handle the adiabatic electronic states for the outer solutions. The inner solutions for the two different values of j do not mix with one another to leading order because of the block diagonal form of (1.14).

One can, of course, take linear combinations of the two solutions described by this theorem. The transition probabilities for Type 2 Avoided Crossings are the same as those presented immediately after the statement of Theorem 3.1 for Type 1 Avoided Crossings.

ACKNOWLEDGMENTS

George Hagedorn wishes to thank the Centre de Physique Théorique. C.N.R.S., Marseille for its hospitality during July 1994. Alain Joye wishes to thank Virginia Polytechnic Institute and State University for its hospitality during August-October 1994 and July-September 1995, and the Fonds National Suisse de la Recherche Scientifique for financial support.

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(Manuscript received May 28, 1996; Revised version received March 10, 1997.)