#### CHAPTER V



#### CONCLUSION

## 5.1 Summary

In Chapter I, we briefly reviewed the one-dimensional white noise model for which the density of states can be computed exactly. This fact was first pointed out by Frisch and Lloyd<sup>3</sup>, and has been subsequently studied in detail by Halperin<sup>2</sup> and Zittartz and Langer.

Defining

$$A(E) = \frac{8}{\pi} (-E),$$
 5.1.1

$$B(E) = \frac{8\sqrt{2}}{3} \cdot \frac{h}{\sqrt{m}} \cdot \frac{(-E)^{3/2}}{\xi},$$
 5.1.2

we can express the exact result for  $\xi \to 0$  as

$$\rho_{as}(E) = \frac{A(E)}{2\xi} \exp\{-\frac{B(E)}{2\xi}\}.$$
 5.1.3

The expression (5.1.3) provides us with an exact expression which can be used to check the validity of more general approximation schemes.

In Chapter II, we examine the method of Halperin and Lax<sup>5,9</sup> used to compute the density of states for the random model. Their theory is based on Schrödinger formulation of quantum mechanics. They assumed that all eigenstates in the low-energy tail of a given energy have the same shape. Using this assumption, they first approximate the density of states<sup>5</sup>, and show that, for their first order approximation,

$$\rho_1(E) = \frac{1}{\sqrt{5}} \cdot \frac{A(E)}{2\xi} \cdot \exp\{-\frac{B(E)}{2\xi}\}$$
 5.1.4

Since their assumption is not exact, the density of states  $\rho_1(E)$  turns out to be too small by a factor of  $\sqrt{5}$  = 2.236. However, in their second order approximation, they calculate the correction C(E), and show that the corrected density of states is given by 9

$$\rho_2(E) = \exp{\{\mu C(E)\}} \rho_1(E),$$

where
 $\mu C(E) = \frac{13}{18}.$ 

5.1.6

It follows that

$$\rho_2(E) = \frac{e^{13/18}}{\sqrt{5}} \cdot \frac{A(E)}{2\xi} \cdot \exp\{-\frac{B(E)}{2\xi}\}, \quad 5.1.7$$

where the unexact factor  $e^{13/18}/\sqrt{5} = 0.921$ .

In Chapter III, we review the Sa-yakanit theory 10,7 and his calculation of the density of states for the random model. His theory is based on Feynman formulation of quantum mechanics. He introduce the non-local harmonic oscillator trial action to model the potential of the random model. He show that for the leading term of the first cumulant, the approximate density of states 7 is given by

$$\rho_{\ell}(E) = \frac{\sqrt{2\pi}}{6} \cdot \frac{A(E)}{2\xi} \cdot \exp^{2(\pi - \frac{\pi}{3})^{1/2}} \cdot \frac{(B(E))}{2\xi}, \quad 5.1.8$$

where the unexact factors are  $\frac{\sqrt{2\pi}}{6}$  = 0.418 and  $(\frac{\pi}{3})^{1/2}$  = 1.0233.

In Chapter IV, we extend Sa-yakanit's theory 7,10 to include the leading term in the second cumulant by using the method of steepest

descent, and show that the corrected density of states corresponding to the correction  $R(\xi t_{o})$  can be put into the form

$$\rho_{e}(E) = \exp\{R(\xi t_{o})\}\rho_{\ell}(E). \qquad 5.1.9$$

For the complete first cumulant, the corrected density of states  $\rho_c(E)$  is  $\rho_1(E)$  and the correction  $R(\xi t_o)$  is  $R_1(\xi t_o)$ , where

$$R_1(\xi t_0) = -\frac{1}{2} + \ln 4;$$
 5.1.10

then  $\rho_1(E) = e^{-1/2} \cdot \frac{4\sqrt{2\pi}}{6} \cdot \frac{A(E)}{2\xi} \cdot \exp\{-(\frac{\pi}{3})^{1/2} \cdot \frac{B(E)}{2\xi}\}.$  5.1.11

For the second cumulant approximation, the corrected density of states  $\rho_c(E)$  is  $\rho_2(E)$ , and the correction  $R(\xi t_o)$  is  $R_1(\xi t_o) + R_2(\xi t_o)$ , where

$$R_{2}(\xi t_{0}) = + \frac{1}{4} \left( \frac{3}{64} + \frac{5}{768} \right) \left( \frac{\pi}{3} \right)^{1/2} \cdot \frac{B(E)}{2\xi} . \qquad 5.1.12$$

$$\rho_{2}(E) = e^{-1/2} \cdot \frac{4\sqrt{2\pi}}{6} \cdot \frac{A(E)}{2\xi} \cdot \exp\left\{ -\left( \frac{3031}{3072} \right) \left( \frac{\pi}{3} \right)^{1/2} \cdot \frac{B(E)}{2\xi} \right\} .$$

5.1.13

and the unexact factors are  $e^{-1/2} \cdot \frac{4\sqrt{2\pi}}{6} = 1.014$  and  $(\frac{3031}{3072})(\frac{\pi}{3})^{1/2} = 1.0097$ .

#### 5.2 Discussion

As seen in the preceding section, all of the approximate density  $^{5,9,7}$  of states do not have the exact numerical factors in front of the exponential term and in the exponent. We have computed the percentage errors of them. Denoting these errors as  $\Delta_1$  and  $\Delta_2$  respectively,

Table 5.1 The percentage errors of the numerical factors in front of the exponential and in the exponent for the correspond density of states: The simbols  $\Delta_1$  and  $\Delta_2$  here denote the percentage errors respectively. For convenience in comparison we write superscript HL to specify the Halperin and Lax results.

p A	` <sup>Δ</sup> 1	Δ <sub>2</sub>	
p <sub>as</sub>	0		
ρHL	55	0	
ρHL 2	-7.9	0	
PR	-58	2.3	
ρ1	1.4	2.3	
P <sub>2</sub>	1.4	0.97	

the errors are listed in Table 5.1

We first consider the Halperin and Lax results : Both p HL and ρ2 have the same zero-values of Δ2's, but have non-zero values for  $\Lambda_1$ 's. The effect of the second order correction changes  $\rho_1^{HL}$  to  $\rho_2^{HL}$ , and changes the values of  $\Delta_1$ 's from -55 to -7.9 %. Next we consider the results of Sa-yakanit's theory : All of the approximate density of states have non-zero values for both  $\Delta_1$ 's and  $\Delta_2$ 's. Both  $\rho_1$  and  $\rho_1$  have the same values of 2.3 % for  $\Delta_2$ 's, but have different values for  $\Delta_1$ 's. The effect of the complete first cumulant changes pe to p, and changes the values of  $\Delta_1$ 's from -58 to + 1.4 %. For  $\rho_2$  the value of  $\Delta_1$  is still 1.4 %, but the value of  $\Delta_2$  is changed from 2.3 to 0.97 % according to the effect of the second cumulant correction. Finally we consider the results from the theories of both Halperin and Lax and Sa-yakanit : We see that our Ap's are not as good as of Halperin and Lax. In principle it is certainly possible to further correct our results until they agree exactly with the exact expression. However, the leading terms in the higher cumulants affect the value of  $\Delta_2$  to a much lesser extent than the leading term in the second cumulant. The values of both  $\Delta_1$  and  $\Delta_2$  of  $\rho_2$  are very close to the exact values. Our result therefore appears to be satisfactory and so we have chosen to stop our calculation at this step. However, if one neglects few percents of  $\Delta_2$ 's of both  $\rho_\ell$  and  $\rho_1$ , one finds that  $\rho_\ell$  and  $\rho_1^{HL}$ ar well as  $\rho_1$  and  $\rho_2$  have nearly the same values of  $\Delta_1$ 's. We see that  $\rho_{\ell}$  and  $\rho_{1}$  are nearly equal to  $\rho_{1}^{HL}$  and  $\rho_{2}^{HL}$  respectively. Our

method 10,7 has, however, several advantages over the method of Halperin and Lax<sup>5,9</sup>.

Firstly, the mathematical details of calculating the Halperin and Lax relation (2.6.4)<sup>9</sup> is very complicated. However, our equivalent relation (4.3.9) is easily obtained by the method of steepest descent.

Secondly, the minimization of the exponent of (2.4.3) in the Halperin and Lax method leads to a nonlinear differential equation

$$-\frac{1}{2}f''(x) - \frac{1}{2}\mu f(x)^3 = Ef(x), \qquad 2.5.5$$

whereas in our work, the minimization leads to a simple algebraic equation

$$E_{\omega} = -\frac{1}{3} E$$
. 3.5.18

We clearly see that (3.5.18) is more simpler than (2.5.5).

Thirdly, the variational equation (3.5.18) which is the minimization of the exponent can be reduced from the maximization of the pressure

$$P(E) = \int_{-\infty}^{E} dE' \int_{-\infty}^{E'} dE'' \rho(E''), \qquad F.1$$

as formulated by Lloyd and Best. This means that we are able to determine the variational parameter  $\omega$  by the Lloyd-Best variational principle 24. No such determination appears in the Halperin and Lax

theory<sup>5,9</sup>. Mathematically our method is more rigorous than their method

Finally, Halperin and Lax's method<sup>5,9</sup> can not be extended to find the density of states for a free electron model, whereas our method can. In the limit of small time, it is obvious that (3.5.5) of our calculation reduces to

$$\rho_{1}(E) = \frac{1}{2\pi h} \left(\frac{m}{2\pi h}\right)^{1/2} \int_{-\infty}^{\infty} dt (it)^{-1/2} \exp \left\{ \frac{iEt}{h} - \frac{1}{2h^{2}} \cdot \frac{\xi}{\sqrt{\pi}L} \cdot t^{2} \right\}.$$
5.2.1

By using the formula (3.5.17) and the asymptotic expansion  $^{23}$  of  $D_p(z)$ 

$$D_{p}(z) = -\frac{\sqrt{2\pi}}{\Gamma(-p)} \cdot e^{p\pi i} e^{z^{2}/4} z^{-p-1},$$
 5.2.2

where I(v) denotes the complete gamma function, we get

$$\rho_1(E) = \frac{1}{2\pi\hbar} \cdot \sqrt{\frac{2m}{E}}$$
 5.2.3

which is usually the density of states for the free electron model.

#### 5.3 Conclusion

As discussed in the preceding section, we see that our method has several advantages over Halperin and Lax's method. However, both of these theories were formulated from the same idea by assuming that the random potential of the system is the same shape. The results of

these two different approaches should not be different. To summerize our results, the curves of Equations (5.1.4),(5.1.7), (5.1.8), (5.1.11), (5.1.13), and (5.1.3) are plotted together in Figure 5.1. For clarity of the curves, Table 5.2 is presented. From this figure and table, we can make the following conclusions:

Firstly, we see that  $\rho_{\ell}$  is not as good as  $\rho_{1}^{HL}$  in comparison with  $\rho_{as}$ . This can be seen easily by considering the percentage errors of  $\rho$  and  $\rho_1^{\rm HL}$ . For example at E = -2, the percentage error of  $\rho$  is -67.5 %, while that of  $\rho_1^{HL}$  is-55.4 %. The reason may be thought of as follows: In our method the shape of potential wells is always fixed in the form of harmonic well, while in Halperin and Lax's method the wave function f is allowed to vary in shape in order to fit the random potential. Fortunately, in one dimension the function f can be solved exactly for the nonlinear differential equation (2.5.5). This fact enable the first order approximation of Halperin and Lax's method to get better results than the leading term of the first cumulant. However, in the complete first cumulant the values of  $\rho_1$  is much better than  $\rho_1^{HL}$ as compared to  $\rho_{as}$ . For example at E = -2, the percentage error of  $\rho_1$ equals -21.2 %, while that of  $\rho_1^{\rm HL}$  equals -67.5 %. This implies that the complete first cumulant give better results than the first order approximation of Halperin and Lax's method. Furthermore, our method can be performed beyond the first cumulant to obtain the leading term of the second cumulant( $\rho_2$ ). By comparing  $\rho_2$  and  $\rho_2^{\rm HL}$  with  $\rho_{as}$ , we find that the results of both  $\rho_2$  and  $\rho_2^{HL}$  are closed to  $\rho_a$ . At high energies, for example at E = -0.5, the percentage error of  $\rho_2$  is extremely small, while that of  $\rho_2^{\rm HL}$  is equal to -8.04 %. At low energies, for example

at E = -2, the percentage error of  $\rho_2^{HL}$  is -8.40 %. From this example, we see that our result  $\rho_2$  is better than  $\rho_2^{HL}$  at high energies, and is nearly the same as  $\rho_2^{HL}$  at low energies. This indicates that our  $\rho_2$  is approximately equal to  $\rho_2^{HL}$ , and that the second cumulant approximation in our method is approximately equal to the second order approximation of Halperin and Lax's method.

Secondly, as discussed in the preceding section, our method has several advantages over Halperin and Lax's method. Firstly, the cumulant correction to the density of states can be calculated directly in our method, but Halperin and Lax must solve the modified Green's function  $\hat{G}$ . Secondly, our method give the variational equation which is easier to solve than that of Halperin and Lax's method. Thirdly, we derive the variational equation from the Lloyd-Best variational principle. Finally, we point out that our method can be extended to high energies to obtain the free-electron density of states. This extension does not appear in Halperin and Lax's method. From all of the above discussion, we conclude that our method is more practical than Halperin and Lax's method. We believe that the idea developed in this thesis based on Sa-yakanit's theory should be applicable to the other problems in three dimensions.

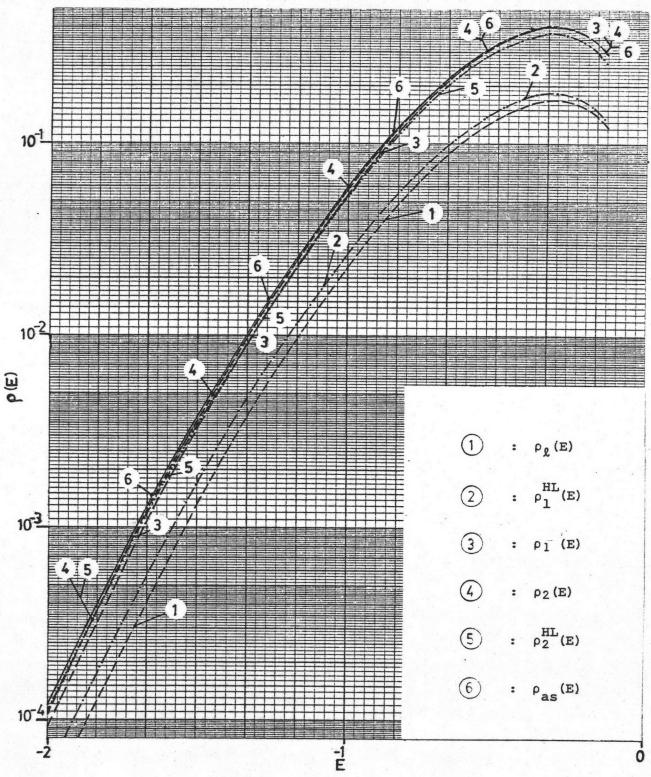


Figure 5.1 Density of states  $\rho(E)$  for the one-dimensional white noise model, compared with its asymptotic form for low energies,  $\rho_{as}(E)$ .

All of the curves are plotted in units of  $\hbar=1$ , m=1, and  $2\xi=1$ . The numbers from one to six indicate the following density of states :

$$\rho_{\mathcal{L}}(E) = \frac{\sqrt{2\pi}}{6} \cdot \frac{A(E)}{2\xi} \cdot \exp\{-(\frac{\pi}{3})^{1/2} \cdot \frac{E(E)}{2\xi}\},$$

$$\rho_{1}^{HL}(E) = \frac{1}{\sqrt{5}} \cdot \frac{A(E)}{2\xi} \exp\{-\frac{E(E)}{2\xi}\},$$

$$\rho_{1}(E) = e^{-1/2} \cdot \frac{4\sqrt{2\pi}}{6} \cdot \frac{A(E)}{2\xi} \cdot \exp\{-(\frac{\pi}{3})^{1/2} \cdot \frac{B(E)}{2\xi}\},$$

$$\rho_{2}(E) = e^{-1/2} \cdot \frac{4\sqrt{2\pi}}{6} \cdot \frac{A(E)}{2\xi} \cdot \exp\{-(\frac{3031}{3072})(\frac{\pi}{3})^{1/2} \cdot \frac{B(E)}{2\xi}\},$$

$$\rho_{2}^{HL}(E) = \frac{e^{13/18}}{\sqrt{5}} \cdot \frac{A(E)}{2\xi} \cdot \exp\{-\frac{E(E)}{2\xi}\},$$
and
$$\rho_{2S}(E) = \frac{A(E)}{2\xi} \cdot \exp\{-\frac{E(E)}{2\xi}\},$$

where the functions A(E) and B(E) are defined by

$$A(E) = \frac{8}{\pi} (-E)$$
,  
 $B(E) = \frac{8\sqrt{2}}{3} (-E)^{3/2}$ 

-E p	ρ <sub>&amp;</sub> (E)	ρ <sup>HL</sup> (E)	ρ <sub>1</sub> (E)	ρ <sub>2</sub> (E)	ρ <sup>HL</sup> (E)	ρ <sub>as</sub> (E)
2.0	3.87-05	5.31-05	9.38-05	1.09-04	1.09-04	1.19-04
1.5	1.33-03	1.67-03	3.23-03	3.55-03	3.45-03	3.74-03
1.0	2.24-02	2.62-02	5.44-02	5.73-02	5.40-02	5.86-02
0.5	1.36-01	1.50-01	3.30-01	3.36-01	3,09-01	3.36-01
1/3	1.69-01	1.84-01	4.09-01	4.13-01	3.78-01	4.11-01
1/6	1.36-01	1.47-01	3.31-01	3.32-01	3.02-01	3.28-01
				,		

Table 5.2 Numerical results of the corresponding density of states which we use to plot Figure 5.1. The forms of all are shown in the page 56. All of the results is written in the computational notation, for example,  $3.87-05 = 3.87 \times 10^{-5}$ .

#### APPENDIX A

# EVALUATION OF $G_0(\vec{x}_2, \vec{x}_1; t)$

To evaluate the non-local harmonic oscillator propagator  $^{20}$   $_{\rm o}$ , we rewrite the trial action  $_{\rm o}$  in the form

$$S_{O}(\omega) = S_{HO} + \frac{m\omega^{2}}{2t} (\int_{0}^{t} d\tau \ \dot{\vec{X}}(\tau))^{2}, \qquad A.1$$

where S is familiarly the simple harmonic oscillator action 12,

$$S_{HO} = \int_{0}^{t} d\tau \frac{m}{2} \{\dot{x}^{2}(\tau) - \omega^{2} \dot{x}(\tau)\}. \qquad A.2$$

The second term on the right-hand side of (A.1) can be converted to an integral form by an identity 23,

$$\exp\{\frac{i}{\hbar}, \frac{m\omega^2}{2t}(\int_0^t d\tau \dot{\vec{X}}(\tau))^2\}$$

$$= \left(\frac{it}{2\pi\hbar m\omega^2}\right)^{d/2} \int_{-\infty}^{\infty} d\vec{f} \exp\left\{-\frac{i}{\hbar}\left(\frac{t}{2m\omega^2} \vec{f}^2 - \int_{0}^{t} d\tau \vec{x}(\tau) \cdot \vec{f}\right)\right\}. A.3$$

Inserting (A.1) and (A.3) into (3.4.3), we then find that the propagator  $G_{0}$  can be expressed as

$$G_{0}(\vec{x}_{2}, \vec{x}_{1}; t) = \left(\frac{it}{2\pi\hbar m\omega^{2}}\right)^{d/2} \cdot \int \mathcal{D}[\vec{x}(\tau)] \int_{-\infty}^{\infty} d\vec{f}.$$

$$\cdot \exp\{\frac{i}{\hbar} \left(S_{HO} - \frac{t}{2m\omega^{2}} \vec{f}^{2} + \int_{0}^{t} d\tau \vec{x}(\tau) \cdot \vec{f}\right)\}. \quad A.4$$

Changing the order of integration, (A.4) becomes

$$G_{O}(\vec{x}_{2},\vec{x}_{1};t) = \left(\frac{it}{2\pi\hbar\omega^{2}}\right)^{d/2} \int_{-\infty}^{\infty} d\vec{r} \exp\left\{-\frac{i}{\hbar} \left(\frac{t}{2m\omega^{2}}\right)\vec{r}^{2}\right\}.$$

$$G_{f}(\vec{x}_{2},\vec{x}_{1};t,\vec{r}), \qquad A.5$$

where

$$G_{\mathbf{f}}(\vec{x}_{2},\vec{x}_{1};t,\vec{f}) = \int \mathcal{D}[\vec{x}(\tau)] \exp\{\frac{i}{\hbar}(S_{HO} + \int_{0}^{t} d\tau \vec{x}(\tau).\vec{f})\}. \quad A.6$$

The propagator (A.6) is the forced harmonic oscillator propagator which is evaluated in the literature 12. With a constant external force  $\vec{f}$  we have

$$G_{\mathbf{f}}(\vec{x}_{2},\vec{x}_{1};t,\vec{f}) = \left(\frac{m\omega}{2\pi i \hbar \sin \omega t}\right)^{d/2} \exp\left\{\frac{i}{\hbar}\left(\frac{m\omega}{4}\left[\cot \frac{\omega t}{2}(\vec{x}_{2}-\vec{x}_{1})^{2}\right]\right) - \tan \frac{\omega t}{2}(\vec{x}_{2}+\vec{x}_{1})^{2}\right\} + \frac{1}{\omega} \cdot \tan \frac{\omega t}{2}(\vec{x}_{2}+\vec{x}_{1}) \cdot \vec{f}$$

$$-\left[\frac{1}{m\omega} \tan \frac{\omega t}{2} - \frac{t}{2m\omega^{2}} \cdot \vec{f}^{2}\right]\right\}. \quad A.7$$

Substituting (A.7) into (A.5), and performing the f - integration. We get

$$G_{O}(\vec{X}_{2},\vec{X}_{1};t) = (\frac{m}{2\pi i \hbar t})^{d/2} \cdot (\frac{\omega t}{2\sin \omega t})^{d}.$$

$$\exp \left[\frac{i}{\hbar} \cdot \frac{m\omega}{4} \cot \frac{\omega t}{2} (\vec{X}_{2} - \vec{X}_{1})^{2}\right]. \quad A.8$$

### APPENDIX B

# EVALUATION OF Sock

The classical action  $S_{\text{o'el}}$  is the extremum of the action  $S_{\text{o}}$ , and the action  $S_{\text{o}}$  can be written as

$$S_{o}^{*}(\vec{x}_{2},\vec{x}_{1};\omega) = S_{o}(\vec{x}_{2},\vec{x}_{1};\omega) + \int_{o}^{t} d\tau \hat{f}(\tau) \cdot \hat{x}(\tau), \qquad B.1$$

where the action  $S_{o}$  corresponds to the classical action  $S_{o,cl}$ , and is the action  $S_{o}(\omega)$  of (3.4.1), and where the external force  $\tilde{f}(\tau)$  is an arbitrary function of time.

To calculate the action  $S_{o,cl}^{*}$ , we substitute (3.4.1) into (B.1), and get

$$S_{o}^{*}(\vec{X}_{2},\vec{X}_{1};\omega) = \int_{0}^{t} d\tau \frac{m}{2} \{\dot{x}^{2}(\tau) - \frac{\omega^{2}}{2t} \cdot \int_{0}^{t} d\sigma [\vec{X}_{c}(\tau) - \vec{X}_{c}(\sigma)]^{2}\}$$

$$+ \int_{0}^{t} d\tau \vec{F}(\tau) \cdot \vec{X}_{c}(\tau).$$
B.2

If  $\vec{X}_{c}(\tau)$  denotes the classical path, then  $S_{o,cl}^{\dagger}$  is given by

$$s_{o,cl}^{\dagger}(\vec{x}_{2},\vec{x}_{1};\omega) = \int_{o}^{t} d\tau \frac{m}{2} (\dot{x}_{c}^{2}(\tau) - \frac{\omega^{2}}{2t} \cdot \int_{o}^{t} d\sigma (\vec{x}_{c}(\tau) - \vec{x}_{c}(\sigma))^{2})$$

$$+ \int_{o}^{t} d\tau \dot{\vec{r}}(\tau) \cdot \dot{\vec{x}}_{c}(\tau).$$
B.3

Taking variational on S', we obtain the classical equation

$$\dot{\bar{\chi}}_{c}(\tau) + \omega^{2}\dot{\bar{\chi}}_{c}(\tau) = \frac{\omega^{2}}{t} \cdot \int_{0}^{t} d\sigma \dot{\bar{\chi}}_{c}(\sigma) + \frac{1}{m}(\tau). \quad B.4$$

If we integrate by parts the first term on the right hand side of (B.3), and use (B.4), we obtain

$$s'_{o,cl}(\vec{x}_2,\vec{x}_1;\omega) = \frac{m}{2} \left[ \vec{x}_c(\tau) \cdot \vec{x}_c(\omega) \right] \Big|_{\tau=0}^{\tau=t} + \frac{1}{2} \int_{0}^{t} d\tau \vec{f}(\tau) \cdot \vec{x}_c(\tau).$$
B.5

The integro-differential equation (B.4) is solvable and so the complete solution of  $\vec{X}_c(\tau)$  with the boundary conditions  $\vec{X}(0) = \vec{X}_1$  and  $\vec{X}(t) = \vec{X}_2$  can be obtained. Inserting  $\vec{X}_c(\tau)$  into (B.5), we get

$$s'_{0,cl}(\vec{x}_{2},\vec{x}_{1};\omega) = \frac{m\omega}{4} \cot \frac{1}{2} |\vec{x}_{2} - \vec{x}_{1}|^{2}$$

$$+ \frac{m\omega}{2\sin\omega t} \left[\frac{2\vec{x}_{2}}{m\omega} \cdot \int_{0}^{t} d\tau \hat{f}(\tau) \left(\sin\omega \tau - 2\sin\omega t \cdot \sin\omega (t-\tau)\sin\omega \tau}{2}\right) + \frac{2\vec{x}_{1}}{m\omega} \cdot \int_{0}^{t} d\tau \hat{f}(\tau) \left(\sin\omega (t-\tau) - 2\sin\omega t \cdot \sin\omega (t-\tau)\sin\omega \tau}{2}\right)$$

$$- \frac{2}{m^{2}\omega^{2}} \cdot \int_{0}^{t} d\tau d\sigma \hat{f}(\tau) \cdot \hat{f}(\sigma) \left(\sin\omega (t-\tau)\sin\omega \sigma\right)$$

$$- 4 \sin\frac{\omega}{2}(t-\tau)\sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega (t-\sigma) \cdot \sin\omega \sigma$$

$$= 4 \sin\frac{\omega}{2}(t-\tau)\sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \sigma$$

$$= 2 \sin\omega t \cdot \sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \sigma$$

$$= 3 \sin\omega t \cdot \sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \sigma$$

$$= 4 \sin\frac{\omega}{2}(t-\tau)\sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \sigma$$

$$= 3 \sin\omega t \cdot \sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \sigma$$

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$$= 4 \sin\frac{\omega}{2}(t-\tau)\sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \sigma$$

$$= 4 \sin\frac{\omega}{2}(t-\tau)\sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \sigma$$

$$= 4 \sin\frac{\omega}{2}(t-\tau)\sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \sigma$$

$$= 4 \sin\frac{\omega}{2}(t-\tau)\sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \tau \cdot \sin\omega \sigma$$

By means of (B.1), the action  $S_0$  becomes the action  $S_0$  when  $\tilde{f}$  equals zero. This means that the classical action  $S_{0,cl}$  becomes the classical action  $S_{0,cl}$  when  $\tilde{f}$ 's equal zero. Hence, we find

$$S_{o,cl}(\vec{x}_2, \vec{x}_1; \omega) = \frac{m\omega}{4} \cot \frac{\omega t}{2} |\vec{x}_2 - \vec{x}_1|^2.$$
 B.7

#### APPENDIX C

# INTEGRAL FORM OF <5">

To find the integral form of  $\langle S^n \rangle$ , we insert (3.3.19) into (3.4.7), and write the average of  $S^n$ , for an integer n, in the form

$$\langle \mathbf{S}^{n} \rangle = \left(\frac{\mathbf{i}}{2h}\right)^{n} \frac{\xi^{n}}{(2\pi)^{dn}} \int_{0}^{t} \int_{0}^{t} \dots \int_{0}^{t} d\tau_{1} d\sigma_{1} \dots d\tau_{n} d\sigma_{n} \int_{\infty}^{\infty} d\vec{k}_{1} \dots \int_{\infty}^{\infty} d\vec{k}_{n} e^{-\frac{L^{2}}{h} \sum_{i=1}^{n} \vec{k}_{i}^{2}} d\vec{k}_{n} e^{-\frac{L^{2}}{h} e^{-\frac{L^{2}}{h} \sum_{i=1}^{n} \vec{k}_{i}^{2}} d\vec{k}_{n} e^{-\frac{L^{2}}{h} e$$

where the quantity  $\xi$  denotes the fluctuation which is defined by (2.5.3), and where the force  $\hat{f}$  is given by

$$\vec{f}(u) \equiv h \sum_{m=1}^{n} \vec{k}_{m} \{\delta(u - \tau_{m}) - \delta(u - \sigma_{m})\}.$$
 C.2

According to Feynman and Hibbs 12, the averaging in (C.1) is called the characteristic functional, and can be done exactly as follows

$$<\exp\{\frac{i}{\hbar}\int_{0}^{\hbar}du\vec{f}(u).\vec{X}(u)\}> = \exp\{\frac{i}{\hbar}(S_{0,cl}^{i}(\vec{X}_{2},\vec{X}_{1};e) - S_{0,cl}(\vec{X}_{2},\vec{X}_{1};e))\}, C.3$$

where the actions So,cl and So,cl are the extremum of the actions So and So,cl respectively, and where the function f depends explicitly on time, and satisfies (B.1).

To compute  $\langle S \rangle$ , we let n = 1 in (C.1), thus, we have

$$\frac{i}{\hbar} < S > = -\frac{1}{2\hbar^2} \cdot \frac{\epsilon}{(2\pi)^d} \cdot \int_0^t \int_0^t d\tau_1 d\sigma_1 \int_{-\infty}^{\infty} d\vec{k}_1 e^{-L^2 \vec{k}^2/4}$$

$$\cdot \langle \exp\{\frac{i}{\hbar} \int_{0}^{t} du \vec{r}(u) \cdot \vec{x}(u)\} \rangle , \qquad C.4$$

where

$$\vec{f}(u) = \hbar \vec{k}_1 \{\delta(u - \tau_1) - \delta(u - \sigma_1)\}. \quad C.5$$

Referring to (B.7) and (B.6), when  $\vec{x}_2 = \vec{x}_1 = 0$ , we can easily show that, for dummy variables  $\tau$  and  $\sigma$ ,

$$S_{0,cl}(0,0,\omega) = 0,$$

$$S_{0,cl}(0,0,\omega) = -\frac{1}{\text{mwsinwt}} \int_{0}^{t} d\tau d\sigma \hat{f}(\tau) \cdot \hat{f}(\sigma) \{\sin \omega(t-\tau) \sin \omega \sigma\}$$

$$-4\sin\underline{\omega}(t-\tau)\sin\underline{\omega}\tau\sin\underline{\omega}(t-\sigma)\sin\underline{\omega}\sigma\}.$$
 C.7

By performing the integration and grouping the trigonometric function, (C.7) becomes

$$\frac{i}{h} s_{o,cl}(o,o;\omega) = -\frac{i\hbar}{m\omega} \cdot \vec{k}_{1}^{2} \cdot \frac{\sin\omega t(|\tau-\sigma|)\sin\omega (t-|\tau-\sigma|)}{2} \cdot C.8$$

or 
$$\frac{1}{h} S_{0,cl}^{*}(0,0;\omega) = \frac{1h}{2m\omega} \cdot \cot \frac{\omega t}{2} k_1^2 \{1 - \Delta(|\tau - \sigma|)\},$$
 C.9

where 
$$\Delta(x) = \tan \omega t \sin \omega x + \cos \omega x$$
. C.10

Substituting (C.8) and (C.6) into (C.3); and substituting (C.3) into (C.4), we find after performing the  $\vec{k}$  - integration

$$\frac{i}{h} < S > = -\frac{1}{2h^2} \cdot \frac{\xi}{(4\pi)^{d/2}} \cdot \int_0^{t} \int_0^{t} d\tau_1 d\sigma_1 \left\{ \frac{L^2}{4} + \frac{i\hbar}{m\omega} \cdot \frac{\sin\omega x}{2} \frac{\sin\omega t}{2} \right\} \cdot C.11$$

C.11

where

$$x = |\tau - \sigma|$$
.

C.12

Letting  $\tau > \sigma$ , and setting

$$\mathbf{x} = \tau - \sigma, \qquad \qquad \mathbf{C.13}$$

and

$$v = \tau + \sigma$$

C.14

the integral in (C.11) can be reduced to

$$\int_{0}^{t} d\tau d\sigma j(x,\omega;t) = t \int_{0}^{t} dx j(x,\omega;t), \qquad C.15$$

where 
$$j(x,\omega;t) = \frac{L^2}{4} + \frac{i\hbar}{m\omega} \cdot \frac{\sin\omega x \cdot \sin\omega(t-x)}{2}$$
. C.16

Finally we have

$$\frac{i}{n} < S > = -\frac{1}{2h^2} \cdot \frac{\xi}{(4\pi)^{d/2}} \cdot t \int_{0}^{t} dx j(x, \omega; t)^{-d/2}.$$
 C.17

One can show that, for n = 2, (C.1) becomes

$$\langle s^2 \rangle = -\frac{1}{4h^2} \cdot \frac{\xi^2}{(4\pi)^d} \cdot \int_{0.00}^{0.00} \int_{0.00}^{0.00} d\tau_1 d\sigma_1 d\tau_2 d\sigma_2 \{\text{Det B}\}^{-d/2}, \quad \text{C.18}$$

where Det B = 
$$\left[ \left\{ \frac{L^2}{4} - \frac{11}{2m\omega} \cdot \cot \frac{\omega t}{2} (1 - \Delta_1) \right\} \left\{ \frac{L^2}{4} - \frac{11}{2m\omega} \cdot \cot \frac{\omega t}{2} (1 - \Delta_2) \right\} - \left( -\frac{11}{2m\omega} \cot \frac{\omega t}{2} \right)^2 \cdot \frac{\Delta_{12}^2}{4} \right],$$
 C.19

$$\Delta_{i} = \Delta(|\tau_{i} - \sigma_{i}|), (i = 1,2),$$
 C.20

and 
$$\Delta_{12} = \Delta_1 + \Delta_2 - \Delta(|\sigma_1 - \tau_2|) - \Delta(|\tau_1 - \sigma_2|).$$

C.21

#### APPENDIX D

#### EVALUATION OF THE INTEGRAL IN <S>

To evaluate the integral in <S> , we define the integral as

$$I(t) \equiv \int_{0}^{1} d\xi \left\{ \sin \underline{wt} \right\} \frac{1}{2} \sin \underline{wt} \left( 1 - \xi \right) / \sin \underline{wt} \right\}. \quad D.1$$

An equivalent form of (D.1) is given by

$$I(t) = (\sin \frac{1}{2}) \int_{0}^{1} d\xi \{\sin^{2} \frac{\omega t}{2} - \sin^{2} \frac{\omega t}{2} (1 - 2\xi)\}.$$
 D.2

Using the transformation

$$\cos \emptyset = \frac{\sin \omega t (1-2\xi)}{\frac{1}{4}}, \qquad D.3$$

the integral (D.2) becomes

$$I(t) = \frac{\frac{1}{\omega t} (\sin \frac{\omega t}{2})^{1/2}}{\int_{0}^{\pi/2} d\phi (1 - \sin^{2} \frac{\omega t}{\mu} \cos^{2} \phi)} \int_{0}^{\pi/2} D.4$$

If we use a relationship between hyperbolic and trigonometric function

$$i sin(ax) = sinh(iax)$$
, D.5

and replace

it = 
$$\beta$$
, D.6

we obtain

$$I(t) = \frac{\frac{1}{4} \left(\frac{\sinh \omega \beta}{2}}{\omega t \sqrt{1} \left(1 + \sinh^2 \omega \beta \right)} K(k), \qquad D.7$$

where 
$$K(k)$$
 
$$\int_{0}^{\pi/2} d\phi (1 - k^2 \sin^2 \phi)^{-1/2}, \quad D.8$$

and 
$$k = \tanh \frac{\omega \beta}{4}$$
. D.9

The complete elliptic integral K(k) in (D.7) can be simplified for computational reason by the Landen's transformation 23

$$K(k) = \frac{1}{(1+k)}K(k^{*}),$$
 D.10

where 
$$k' = \frac{2\sqrt{k}}{1+k} = \left(\frac{2}{1+\coth\omega\beta}\right)^{1/2}$$
. D.11

In the limit of large time, we have

$$I(t) = \frac{1}{\omega t \sqrt{2i}} \cdot k' K(k'), \qquad D.12$$

where the limiting value of the elliptic integral K(k') is given by 27

$$\lim_{\mathbf{k}^{\circ} \to 1} K(\mathbf{k}^{\circ}) = \frac{1}{2} \ln\left(\frac{16}{\mathbf{k}^{2}}\right), \qquad D.13$$

and here 
$$(k^i)^2 = 1 - k^2 = 1 - e^{-\omega \beta}$$
 D.14

Finally we find

$$I(t) = \left(\frac{2}{\omega t}\right) \left(\frac{2}{i}\right)^{1/2} \left(\frac{i\omega t}{2} + \ln 4\right) \qquad D.15$$

#### APPENDIX E

### EVALUATION OF THE A-INTEGRAL

The  $\Delta$ -integral occurs when we consider the leading terms in (4.5.17), and is conveniently defined for an integer **a** by

$$I_{n} \equiv \int_{0}^{1} \int_{0}^{1} dt_{1} ds_{1} \Delta^{n}(|t_{1}-s_{1}|), \qquad E.1$$

where  $\Delta$  is related to (C.10) by the relation; x = yt. The integral (E.1) can be simplified by an elementary conformal mapping, which is similar to (C.13) and (C.14), into the form

$$I_n = \int_0^1 d\xi \, \Delta^n(\xi) , \qquad E.2$$

where 
$$\Delta(\xi) = (e^{-i\omega t}\xi + e^{-i\omega t(1-\xi)})/(1 + e^{-i\omega t})$$
. E.3

In the limit of large time, we can neglect  $^{12}$  the exponential  $e^{-i\omega t}$  and the integral (E.2), for n = 1, is

$$I_1 = \int_0^1 d\xi \ \Delta(\xi) = \frac{2}{i\omega t}.$$
 E.4

Similarly It can be proved in general that

$$I_{n} = \left(\frac{2}{i\omega t}\right)\left(\frac{1}{n}\right).$$
 E.5

Next we shall show that the leading terms in (4.5.17) are just the  $I_{2n}$ 's. To show this, we first consider an integral

$$J_{2n} = \iiint \int \int \int \int \Delta_{12}^{2n} dt_1 ds_1 dt_2 ds_2,$$
 E.6

where  $\Delta_{12}$  is the same as  $\Delta_{12}$  in (4.5.17). If one symbolizes the  $\Delta$  by a line between its points<sup>8</sup>, for example,

$$\Delta(|t_1-s_1|) = \begin{bmatrix} s_1 \\ t_1 \end{bmatrix} \cdot s_2 \\ t_2 \end{bmatrix}$$
, E.7

while a product of them can be expressed, e.g.,

$$\Delta^{n}(|t_{1}-s_{1}|)\Delta^{m}(|t_{1}-s_{2}|) = n m$$
. E.8

where each of n and m is an integer.

Returning to (E.6), we now consider the orders of the integrals

$$\iiint \left(\frac{n}{2}\right) dt_1 ds_1 dt_2 ds_2 = \iiint \left(\frac{n}{2}\right) dt_1 dt_2 \sim t^{-1}, E.9$$

$$\iiint \left(\frac{n}{2}\right) dt_1 ds_1 dt_2 ds_2 \sim \left\{\iiint \left(\frac{n}{2}\right) dt_1 dt_2\right\}^2 \sim t^{-2}, E.10$$

$$\iiint \left(\frac{n}{2}\right) dt_1 ds_1 dt_2 ds_2 \sim t^{-2}.$$

$$E.11$$

We emphasize here that the integrals which yield the results of the order of t<sup>-1</sup> come only from the integral of a single line type shown in(E.9). According to (E.8), we can write (C.21) in the symbolized form as

$$\Delta_{12}^{2} = \left( \underbrace{2} + \underbrace{2} + \underbrace{2} + \underbrace{2} + \underbrace{2} \right) + 2 \left( \underbrace{-} + \underbrace{\times} \right)$$

$$-2 \left( \underbrace{/} + \underbrace{-} + \underbrace{/} + \underbrace{-} + \underbrace{-} \right). \qquad E.11$$

Because of the symmetry of J2n, (E.11) has an equivalent form

$$\Delta_{12}^{2} = 4\left(\underline{2}\right) + 4\left(\underline{2}\right) - 8\left(\underline{2}\right). \qquad E.12$$

The expression (E.12) implies that the leading term of  $J_2$  is of the order of  $t^{-1}$ . Similarly one can prove in general that the leading term of  $J_{2n}$  is also of the order of  $t^{-1}$ . It means that in the leading term approximation we are look at  $I_{2n}$  instead of  $J_{2n}$ ; thus, we have

$$J_{2n} \rightarrow 4I_{2n}$$
 . E.13

#### APPENDIX F

# LLOYD-BEST VARIATIONAL PRINCIPLE

Lloyd and Best  $^{24}$  formulate an exact variational principle for calculating the density of states over an entire energy range. They showed that it is not  $\rho(E)$  or  $\ln \rho(E)$  which must be maximized, but the functional

$$P(E) = \int_{-\infty}^{E} dE' \int_{-\infty}^{E'} dE'' \rho(E'').$$
 F.1

The functional (F.1) is a pressure of a hypothetical free-fermion system, and has a fascinating analytic relationship to Boltzmann's H-functional.

Recently Sa-yakanit showed in his work that the pressure (F.1) can be simplified by taking an integration by parts

$$P(E) = E' \int_{-\infty}^{E'} dE'' \rho(E'') \Big|_{E'=-\infty}^{E'=E} \int_{-\infty}^{E} dE' E' \rho(E'), \qquad F.2$$

and using the boundary condition that

$$\rho(-\infty) = 0.$$
 F.3

It then follows that

$$P(E,\omega) = \int_{-\infty}^{E} dE'(E - E')\rho(E',\omega). \qquad F.4$$

The result (F.4) is obtained by replacing P(E) and  $\rho$ (E) by P(E, $\omega$ ) and  $\rho$ (E, $\omega$ ) respectively for special emphasis on the variational parameter

According to Section 4.2, if we formally set

$$\frac{\alpha^2}{4\beta\xi} = \frac{x^2}{p}, \qquad F.5$$

where 
$$x = \frac{E_{\omega}}{h} - E$$
, F.6

and 
$$p \sim \omega^{1/2}$$
, F.7

then we get

$$\rho(x',\omega) \sim \omega^{1/2}(x')^{1/2} \exp\{-(x')^2/p\}$$
. F.8

By means of (F.6), we now write (F.4) as

$$P(x',\omega) = \int_{x}^{\infty} dx'(x'-x)\rho(x',\omega). \qquad F.9$$

To perform the variational principle, we use the transformation

$$(x')^2 = pz', F.10$$

and then the maximization of the pressure leads us to

$$0 = \frac{d}{d\omega} \omega^{9/8} \int_{z}^{\infty} dz' e^{-z'} \{(z')^{1/4} - (z')^{-1/4}, z^{1/2}\}.$$
 F.11

Each integral in (F.11) is an incomplete gamma function 27 which is defined by

$$\Gamma(\alpha, y) \equiv \int_{y}^{\infty} dt e^{-t} t^{\alpha-1}$$
. F.12

By using (F.12), (F.11) becomes

$$0 = \frac{d}{dm} \omega^{9/8} \{ \Gamma(\frac{5}{4}, z) - z^{1/2} \Gamma(\frac{3}{4}, z) \}.$$
 F.13

If one performs the differentiation in (F.13), and uses the identity 27

$$\frac{d}{dy} \left\{ \Gamma(\alpha, y) \right\} = -y^{\alpha-1} e^{-y}, \qquad F.14$$

then

$$0 = \left[\frac{9}{4} z^{1/2} \left\{ \Gamma(\frac{5}{4}, z) - z^{1/2} \Gamma(\frac{3}{4}, z) \right\} - \omega \Gamma(\frac{3}{4}, z) \cdot \frac{\dot{c}}{d\omega} \right\}. \quad F.15$$

Using the asymptotic expansion of the incomplete gamma fun:tion 23

$$\Gamma(\alpha,y) \simeq y^{\alpha-1} e^{-y}$$
, F.16  
 $y \to +\infty$ 

one can show that in the case z > 000

$$\Gamma(\frac{5}{h}, z) - z^{1/2} \Gamma(\frac{3}{h}, z) = 0.$$
 F.17

Returning to (F.5) and (F.10), we find that z is the exponent of the exponential in the density of states

$$z = \frac{\alpha^2}{4\beta\xi} .$$
 F.]8

With (F.18) in the limit  $\xi \to 0$ , and using (F.17); (F.15) reduces to the familiar minimization of the exponent

$$\frac{dz}{d\omega} = \frac{d}{d\omega} \left( \frac{\alpha^2}{4\beta \xi} \right) = 0.$$
 F.19

We thus get the usual variational equation

$$E_{\omega} = -\frac{1}{3} E$$
. F.20

Next we consider the discussion on the end of Section 4.3. When the corrections are included in the density of states, we find the Lloyd-Best variational principle 24 by expanding the exponential of the corrections in series and going beyond the previous steps. This procedure leads us to

$$\rho_{n}(x',\omega) \sim F_{n}(\omega)(x')^{n}(x')^{1/2} \exp\{-(x')^{2}/p\}, F.21$$

where  $\rho_n(x',\omega)$  is the n th term of the density of states in power of x', and  $F_n(\omega)$  is the function of  $\omega$  corresponding to  $\rho_n(x',\omega)$ . Using pressure  $P_n(x',\omega)$  defined for (F.21), and we find  $P_n'(x',\omega)$  as

$$\frac{d}{d\omega} P_{n}(x',\omega) = \frac{d}{d\omega} E_{n}(\omega) \{ \Gamma(\frac{5+2n}{4},z) - z^{1/2} \Gamma(\frac{3+2n}{4},z) \}$$

$$-E_{n}(\omega) \{ \frac{1}{2\sqrt{z}} \Gamma(\frac{3+2n}{4},z) \} \frac{dz}{d\omega} , \qquad F.22$$

where  $E_n(\omega)$  is the function of  $\omega$  corresponding to  $F_n(\omega)$ . According to (F.16), we see that the two terms in the first parenthesis {} in (F.22) cancel each other in the limit  $z \to \infty$ , or equivalently  $\xi \to \infty$ . Moreover, we find the maximization of the pressure  $P(x',\omega)$  is in the form

$$0 = \frac{d}{d\omega} P(x', \omega) = \sum_{n=1}^{\infty} C_n \frac{d}{d\omega} P_n(x', \omega), \qquad F.23$$

where  $C_n$  is the constant of (F.21), and is explicitly independent of  $\omega$ . In the limit of small  $\xi$ , (F.22) implies that (F.23) has a common factor  $\frac{dz}{d\omega}$ . It means that (F.19) is held, and (F.20) is also true.

Finally we should emphasize that (F.21) and (F.22) are also true even if the n in  $(x')^n$  and  $\Gamma(\alpha,y)$  is a fraction. This occurs when we include the effect of changing dimension of the white noise model in n. It implies that our proof is valid for all dimensionalities of the white noise model.