



Development of a QMC code to tackle interacting electronic systems in 2D with application to TMD nanoribbons

Francisco Monteiro de Oliveira Brito

Thesis to obtain the Master of Science Degree in

Physics Engineering

Supervisor(s): Prof. Eduardo Filipe Vieira de Castro Prof. João Manuel Viana Parente Lopes

Examination Committee

Chairperson: Prof. Lorem
Supervisor: Prof. Eduardo Filipe Vieira de Castro
Co-Supervisor: Prof. João Manuel Viana Parente Lopes
Members of the Committee: Dr. Lorem Ipsum
Prof. Lorem Ipsum



Acknowledgments

I would like to thank the Academy, bla bla bla..

Abstract

The Objective of this Work ... (English)

Keywords

Keywords (English)

Resumo

O objectivo deste trabalho ... (Português)

Palavras Chave

Palavras-Chave (Português)

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Abbreviations

acro Dummy Acronym

List of Symbols

1

Introduction

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1.1 Motivation

Electron correlations are so cool!

1.2 State of The Art

State of The Art Section.

1.2.1 Dummy Subsection A

State of Art Subsection A

1.2.2 Dummy Subsection B

State of Art Subsection B

1.3 Original Contributions

Contributions Section.

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Outline Section.

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A Chapter

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2.1 Section A

2.1.1 Subsection A

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acro

Dummy Acronym

acros

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2.1.2 Subsection B



Figure 2.1: Dummy Figure Caption.

Remember you can change the reference style. Another dummy citation [?].

2.2 Section B

2.2.1 Subsection A

The model described can also be represented as

$$\dot{\mathbf{x}}(t) = \mathbf{T}\mathbf{z}(y), \ \mathbf{y}(0) = \mathbf{y}_0, \ z \ge 0$$
(2.1)

where

$$\mathbf{A} = \begin{bmatrix} -(a_{12} + a_{10}) & a_{21} \\ a_{12} & -(a_{21} + a_{20}) \end{bmatrix}, \ \mathbf{x} = \begin{bmatrix} x_1 \\ x_2 \end{bmatrix}$$
 (2.2)

2.2.2 Subsection B

Table 2.1: Dummy Table.

Vendor Name	Short Name	Commercial Name	Manufacturer
	ABC	ABC [®]	ABC SA
Text in Multiple Row	DEF	DEF®	DEF SA
	GHF	GHF®	GHF SA
Text in Single Row	IJK	IJK®	IJK SA
Frescos SA	LMN	LMN®	LMN SA
Carros Lda.	Text in Multiple Column		



Finite Temperature Auxiliary Field Quantum Monte Carlo

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3.1 Section A

3.1.1 Subsection A

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3.1.2 Auxiliary field QMC

We will now discuss a numerical method to simulate the Hubbard model. Among the various methods belonging to the family of QMC methods, auxiliary field QMC allows us to circumvent the sign problem for the half filled Hubbard model. The sign problem is an uncontrolled numerical error due to the antisymmetry of the many-electron wave function, leading to oscillations in the sign of the quantities that we are interested in measuring. These oscillations deem the algorithm exponentially complex in the size of the system, in general, but it possible to overcome this hurdle for a class of models, namely the Hubbard model at half filling. The difficulty lies in computing averages of quantities that are very close to zero, on average, but have a large variance, i.e. $\sigma_X/\langle X\rangle\gg 1$.

We seek a computable approximation of the projection operator \mathcal{P} defined in equation (??). As we shall see, it is found by using a discrete Hubbard-Stratonovich transformation. This transformation introduces an auxiliary field (consisting basically of Ising spins), and we use Monte Carlo to sample configurations from the distribution corresponding to this *classical* configuration space.

For now, let us assume half filling $\mu=0$, so that there is no sign problem. In fact, many interesting phenomena occur at half filling, for example magnetic ordering and the Mott metal-insulator transition.

3.1.2.A Hubbard-Stratonovich transformation

A – In section **??**, we found exact solutions for particular instances of the Hubbard model by finding a closed form for the partition function [**?**]. When devising a numerical method, a good sanity check is to verify that it satisfactorily approximates the partition function.

The operators \mathcal{H}_K and \mathcal{H}_V of equation (??) do not commute. This impedes us from factorizing the exponential of their sum $e^{-\beta(\mathcal{H}_K+\mathcal{H}_V)}$ exactly. The Trotter-Suzuki decomposition leads to the sought approximate factorization that is used to approximate the partition function. Quantum states evolve according to

$$|\psi(\tau)\rangle = e^{-\tau \mathcal{H}} |\psi(0)\rangle, \tag{3.1}$$

where $\tau=it$ is the imaginary time. Recall that Diffusion Monte Carlo is based on this imaginary time evolution, filtering out the ground state as the state that takes longer to vanish exponentially. We now seek a finite temperature method. To find it, we invoke an analogy with the evolution of a quantum system according to the previous equation.

Taking the scalar product with a position eigenstate $\langle \boldsymbol{x}|$, we obtain $\psi(\boldsymbol{x},\tau) = \langle \boldsymbol{x}|\psi(\tau)\rangle$. Using the closure relation $\int d\boldsymbol{y} \, |\boldsymbol{y}\rangle \, \langle \boldsymbol{y}| = 1$, we get

$$\psi(\boldsymbol{x},\tau) = \int d\boldsymbol{y} \left\langle \boldsymbol{x} | e^{-\tau \mathcal{H}} | \boldsymbol{y} \right\rangle \psi(\boldsymbol{y},0)$$
(3.2)

The wave function at position \boldsymbol{x} and time t may be obtained by this equation as long as we know the wave function at $\tau=0,\,\psi(\boldsymbol{y},0)$ for all points in space \boldsymbol{y} . The evolution operator matrix element, or Green function,

$$G(\boldsymbol{x}, \tau | \boldsymbol{y}, 0) \equiv \langle \boldsymbol{x} | e^{-\tau \mathcal{H}} | \boldsymbol{y} \rangle,$$
 (3.3)

as the wave function, satisfies the Schrödinger equation, with the initial condition $\psi(y,0) = \delta(x-y)$. It is then the probability of presence at x,t of a wave packet centered at y at t=0. Note that the solution of the Schrödinger equation is then analogous to that of a diffusion equation (that in turn one may obtain as the continuum limit of a random walk). We may write G as a linear combination of the eigenstates of the Hamiltonian

$$G(\boldsymbol{x}, \tau | \boldsymbol{y}, 0) = \sum_{\alpha} \psi_{\alpha}^{\star}(\boldsymbol{y}) \psi_{\alpha}(\boldsymbol{x}) e^{-E_{\alpha}\tau},$$
(3.4)

where we immediately note a striking similarity with equation (??). The correspondence $\psi(x,\tau) \mapsto Z_{\beta}$, where $\tau \mapsto \beta$, with respect to section ?? makes the analogy evident. x has no correspondence because it is not a parameter, it is just an arbitrary position that we fixed for the sake of the argument.

Computing the partition function at finite temperature

$$Z_{\beta} = \text{Tr}(e^{-\beta \mathcal{H}}) \tag{3.5}$$

is analogous to computing the Green function of a quantum system evolving in imaginary time. The inverse temperature β now represents the imaginary time $\tau=it$, and Z_{β} may be simply thought of as the wave function of the analogous quantum system at imaginary time (temperature) β .

This expression is not very amenable to numerical computation since it contains an exponential of a sum of operators $\mathcal{H}_K + \mathcal{H}_V$, which is not factorizable and involves computing an infinite number of commutators containing these two operators, as per the Zassenhaus formula, valid for any two generic operators X and Y:

$$e^{\delta(X+Y)} = e^{\delta X} e^{\delta Y} e^{-\frac{\delta^2}{2}[X,Y]} e^{\frac{\delta^3}{6}(2[Y,[X,Y]]+[X,[X,Y]])} e^{-\frac{\delta^4}{24}([[[X,Y],X],X]+3[[[X,Y],X],Y]+3[[[X,Y],Y],Y])}$$
(3.6)

where $\delta \in \mathbb{C}$ is an expansion parameter.

Dividing the imaginary time interval $[0, \beta]$ into L equal sub-intervals of width $\Delta \tau = \beta/L$, we obtain

$$Z = \text{Tr}\bigg(\prod_{l=1}^{L} e^{-\Delta \tau \mathcal{H}}\bigg),\tag{3.7}$$

which is now a product of exponentials of operators multiplied by a constant that can be made small by increasing L. The Trotter-Suzuki decomposition follows from truncating equation (3.6), and keeping only the first order term in t, i.e. the one in $\Delta \tau$ in our case.

$$Z = \text{Tr}\left(\prod_{l=1}^{L} e^{-\Delta \tau \mathcal{H}_K} e^{-\Delta \tau \mathcal{H}_V}\right) + \mathcal{O}(\Delta \tau^2)$$
(3.8)

The kinetic energy term is quadratic in the fermion operators, and is spin-independent and thus may be separated into spin up and spin down components

$$e^{-\Delta \tau \mathcal{H}_K} = e^{-\Delta \tau \mathcal{H}_{K\uparrow}} e^{-\Delta \tau \mathcal{H}_{K\downarrow}}, \tag{3.9}$$

where $\mathcal{H}_{K_{\sigma}}=-toldsymbol{c}_{\sigma}^{\dagger}oldsymbol{K}oldsymbol{c}_{\sigma}$.

The potential energy term, however, is quartic. Surprisingly, it is possible to express it in quadratic form by introducing an extra degree of freedom, the so called *Hubbard-Stratonovich (HS) field* $h \equiv (h_i)_{i=1}^N$, in which each element is essentially an Ising spin. First, note that number operators on different sites commute, so that we have

$$e^{-\Delta \tau \mathcal{H}_{V}} = e^{-U\Delta \tau \sum_{i=1}^{N} (n_{i\uparrow} - 1/2)(n_{i\downarrow} - 1/2)}$$

$$= \prod_{i} e^{-U\Delta \tau (n_{i\uparrow} - 1/2)(n_{i\downarrow} - 1/2)}$$
(3.10)

Now we introduce the discrete Hubbard Stratonovich transformation for U>0 that allows us to recast the equation above in terms of a non-interacting quadratic term $(n_{i\uparrow}-n_{i\downarrow})$.

$$e^{-U\Delta\tau(n_{i\uparrow}-1/2)(n_{i\downarrow}-1/2)} = c_U \sum_{h_i=\pm 1} e^{\nu h_i(n_{i\uparrow}-n_{i\downarrow})},$$
(3.11)

where $c_U=\frac{1}{2}e^{-\frac{U\Delta au}{4}}$ and $u=\mathrm{arcosh}(e^{\frac{U\Delta au}{2}}).$

To prove this identity, let us write down how the operators $(n_{i\uparrow}-1/2)(n_{i\downarrow}-1/2)$ and $(n_{i\uparrow}-n_{i\downarrow})$ act on a state on a given site.

$$(n_{i\uparrow} - 1/2)(n_{i\downarrow} - 1/2) \begin{cases} | \rangle = \frac{1}{4} | \rangle \\ | \uparrow \rangle = -\frac{1}{4} | \uparrow \rangle \\ | \downarrow \rangle = -\frac{1}{4} | \downarrow \rangle \\ | \uparrow \downarrow \rangle = \frac{1}{4} | \uparrow \downarrow \rangle \end{cases}$$

$$(n_{i\uparrow} - n_{i\downarrow}) \begin{cases} | \rangle = 0 | \rangle \\ | \uparrow \rangle = | \uparrow \rangle \\ | \downarrow \rangle = | \downarrow \rangle \\ | \uparrow \downarrow \rangle = 0 | \uparrow \downarrow \rangle \end{cases}$$

$$(3.12)$$

Now we simply compare the action of the operators on the left hand side and on the right hand side of equation (3.11) and find the desired relation by defining

$$\cosh \nu = \frac{e^{\nu} + e^{-\nu}}{2} \equiv e^{\frac{U\Delta\tau}{2}} \tag{3.13}$$

$$e^{-U\Delta\tau(n_{i\uparrow}-1/2)(n_{i\downarrow}-1/2)} |\psi\rangle = e^{-\frac{U\Delta\tau}{4}} |\psi\rangle , |\psi\rangle = |\rangle, |\uparrow\downarrow\rangle$$

$$e^{-U\Delta\tau(n_{i\uparrow}-1/2)(n_{i\downarrow}-1/2)} |\uparrow(\downarrow)\rangle = e^{\frac{U\Delta\tau}{4}} |\uparrow(\downarrow)\rangle$$

$$c_{U} \sum_{h_{i}=\pm 1} e^{\nu h_{i}(n_{i\uparrow}-n_{i\downarrow})} |\psi\rangle = e^{-\frac{U\Delta\tau}{4}} |\psi\rangle , |\psi\rangle = |\rangle, |\uparrow\downarrow\rangle$$

$$c_{U} \sum_{h_{i}=\pm 1} e^{\nu h_{i}(n_{i\uparrow}-n_{i\downarrow})} |\uparrow(\downarrow)\rangle = \frac{e^{\nu}+e^{-\nu}}{2} e^{-\frac{U\Delta\tau}{4}} |\uparrow(\downarrow)\rangle$$
(3.14)

Note that we require U>0 so that there exists $\nu\in\mathbb{R}$ such that $\cosh\nu=e^{U\Delta\tau/2}$. A similar reasoning could be made for U<0. Additionally, other transformations that recast other types of quartic terms in terms of quadratic ones exist, but we shall not need them in what follows [?]. The transformation we derived is the one we will use throughout.

We have now made progress. At the expense of introducing an extra N-dimensional HS-field h, we obtained an *exact* representation of the quartic term in terms of quadratic terms [?].

$$e^{-\Delta \tau \mathcal{H}_V} = \prod_{i=1}^N \left(c_U \sum_{h_i = \pm 1} e^{\nu h_i (n_{i\uparrow} - n_{i\downarrow})} \right), \tag{3.15}$$

which can be manipulated to arrive at a more compact form.

$$e^{-\Delta\tau\mathcal{H}_{V}} = (c_{U})^{N} \sum_{h_{i}=\pm 1} e^{\nu h_{i}(n_{1\uparrow}-n_{1\downarrow})} \sum_{h_{i}=\pm 1} e^{\nu h_{i}(n_{2\uparrow}-n_{2\downarrow})}$$

$$\dots \sum_{h_{i}=\pm 1} e^{\nu h_{i}(n_{N\uparrow}-n_{N\downarrow})}$$

$$= (c_{U})^{N} \sum_{h_{i}=\pm 1} e^{\sum_{i=1}^{N} [(\nu h_{i}(n_{i\uparrow}-n_{i\downarrow})]}$$

$$\equiv (c_{U})^{N} \operatorname{Tr}_{h} e^{\sum_{i=1}^{N} [(\nu h_{i}(n_{i\uparrow}-n_{i\downarrow})]}$$

$$= (c_{U})^{N} \operatorname{Tr}_{h} e^{\sum_{i=1}^{N} \nu h_{i} n_{i\uparrow}} e^{-\sum_{i=1}^{N} \nu h_{i} n_{i\uparrow}}$$

$$= (c_{U})^{N} \operatorname{Tr}_{h} (e^{\mathcal{H}_{V\uparrow}} e^{\mathcal{H}_{V\downarrow}}),$$

$$(3.16)$$

where the spin up and spin down operators \mathcal{H}_{V_σ} are defined as follows

$$\mathcal{H}_{V\sigma} = \sum_{i=1}^{N} \nu h_i n_{i\sigma} = \sigma \nu c_{\sigma}^{\dagger} V(\mathbf{h}) c_{\sigma}, \qquad (3.17)$$

with V(h) being simply the HS-field put into a diagonal $N \times N$ matrix: $V(h) \equiv \text{diag}(h_1, h_2, ..., h_N)$.

For each imaginary time slice l (where $l \in [1, L]$) we may define a HS-field h_l , which in turn specifies V_l and $\mathcal{H}^l_{V_\sigma}$. We may now replace the result of equation (3.16) in equation (3.8), and exchange the traces to obtain

$$Z_{h} = (c_{U})^{NL} \operatorname{Tr}_{h} \operatorname{Tr} \left[\prod_{l=1}^{L} \underbrace{\left(e^{-\Delta \tau \mathcal{H}_{K_{\uparrow}}} e^{\mathcal{H}_{V_{\uparrow}}^{l}} \right)}_{B_{l,\uparrow}(h_{l})} \right]$$

$$\underbrace{\left(e^{-\Delta \tau \mathcal{H}_{K_{\downarrow}}} e^{\mathcal{H}_{V_{\downarrow}}^{l}} \right)}_{B_{l,\uparrow}(h_{l})} \right],$$
(3.18)

where all operators are now quadratic in the fermion operators:

$$\mathcal{H}_{K_{\sigma}} = -t \boldsymbol{c}_{\sigma}^{\dagger} \boldsymbol{K} \boldsymbol{c}_{\sigma}$$

$$\mathcal{H}_{V}^{l} = \sigma \nu \boldsymbol{c}_{\sigma}^{\dagger} \boldsymbol{V}_{l}(\boldsymbol{h}_{l}) \boldsymbol{c}_{\sigma}$$
(3.19)

for $\sigma = \pm 1$ and $V_l(h_l) = \text{diag}(h_{l,1}, h_{l,2}, ..., h_{l,N})$.

Furthermore, we have defined the B-matrices

$$\boldsymbol{B}_{l,\sigma}(\boldsymbol{h}_l) = e^{t\Delta\tau \boldsymbol{K}} e^{\sigma\nu \boldsymbol{V}_l(\boldsymbol{h}_l)}$$
(3.20)

Note that the argument of the first exponential is positive since K is defined so that its entries are 0's and 1's; otherwise (defining K with 0's and -1's) it would be negative.

The problem of computing the partition has been reduced to computing the trace of a product of exponentials of quadratic forms. Thus, we may still rewrite equation (3.18) by making use of the following identity.

Let \mathcal{H}_l be quadratic forms of the fermion operators:

$$\mathcal{H}_l = c_i^{\dagger}(H_l)_{ij}c_j,\tag{3.21}$$

where the summation is implied, and where H_l are real matrices. Then, the following identity holds

$$Tr[e^{-\mathcal{H}_1}e^{-\mathcal{H}_2}...e^{-\mathcal{H}_L}] = \det(I + e^{-H_L}e^{-H_{L-1}}...e^{-H_1})$$
(3.22)

For simplicity, we present the proof for a simpler case, corresponding to a single B-matrix, i.e. a product of exponentials of two quadratic operators [?]. It could then be easily extended to the more general case. Let the two arbitrary real matrices be M and N. Then, a particular case of the previous identity is

$$\operatorname{Tr}\left[e^{-c_i^{\dagger}M_{ij}c_j}e^{-c_i^{\dagger}N_{ij}c_j}\right] = \det(\boldsymbol{I} + e^{-\boldsymbol{M}}e^{-\boldsymbol{N}}), \tag{3.23}$$

where a summation over repeated indices is implied, as it wil be throughout this proof.

To prove this identity, we start by proving that

$$e^{-c_i^{\dagger} M_{ij} c_j} e^{-c_i^{\dagger} N_{ij} c_j} = e^{-\sum_{\nu} c_{\nu}^{\dagger} \rho_{\nu} c_{\nu}}, \tag{3.24}$$

where $\lambda_{\nu}=e^{-\rho_{\nu}}$ are the eigenvalues of the matrix $e^{-M}e^{-N}$.

The proof consists of showing that any many-particle state are propagated in the same way when acted upon by any of these two operators, i.e. the LHS operator leads the system to the same state as the RHS operator.

A generic single-particle state reads

$$|\phi\rangle = \sum_{i} a_{j} c_{j}^{\dagger} |0\rangle , \qquad (3.25)$$

where a_j are arbitrary coefficients, and $|0\rangle$ is the vacuum state.

Let $\{|\mu\rangle\}$ be the basis in which the matrix N is diagonal. Using Dirac notation, we then have

$$N = \sum_{\mu} |\mu\rangle \, n_{\mu} \, \langle \mu| \tag{3.26}$$

Defining new fermionic operators

$$c_{\mu} = \sum_{j} \langle \mu | j \rangle c_{j}$$

$$c_{\mu}^{\dagger} = \sum_{j} \langle j | \mu \rangle c_{j}^{\dagger},$$
(3.27)

which may be inverted to obtain

$$c_{j} = \sum_{\mu} \langle j | \mu \rangle c_{\mu}$$

$$c_{j}^{\dagger} = \sum_{\mu} \langle \mu | j \rangle c_{\mu}^{\dagger},$$
(3.28)

Now we prove yet another identity that goes into proving equation (3.24).

$$e^{-c_i^{\dagger} N_{ij} c_j} = \prod_{\mu} \left[\mathbb{1} + (e^{-n_{\mu}} - 1) c_{\mu}^{\dagger} c_{\mu} \right]$$
 (3.29)

$$\begin{split} &\exp(-c_i^\dagger N_{ij} c_j) = \exp(-\sum_{\mu\nu} \langle \mu | i \rangle \, c_\mu^\dagger N_{ij} \, \langle j | \nu \rangle \, c_\nu) \\ &= \exp(-\sum_{ij} \sum_{\mu\nu\sigma} \langle \mu | i \rangle \, \langle i | \sigma \rangle \, c_\mu^\dagger n_\sigma \, \langle \sigma | j \rangle \, \langle j | \nu \rangle \, c_\nu), \\ &(\text{using the closure relation} \sum_i | i \rangle \, \langle i | = \mathbb{1}) \\ &= \exp(-\sum_{\mu\nu\sigma} \overbrace{\langle \mu | \sigma \rangle}^{\delta_{\mu\sigma}} \, c_\mu^\dagger n_\sigma \, \overbrace{\langle \sigma | \nu \rangle}^{\delta_{\sigma\nu}} \, c_\nu) \\ &= \exp(-\sum_\mu c_\mu^\dagger n_\mu c_\mu) \\ &= \prod_\mu e^{-n_\mu \hat{n}_\mu} \\ &= \prod_\mu [\mathbb{1} + (-n_\mu \hat{n}_\mu + \frac{n_\mu^2}{2!} \hat{n}_\mu^2 - \frac{n_\mu^3}{3!} \hat{n}_\mu^3 + \ldots] \\ &= \prod [\mathbb{1} + (-n_\mu + \frac{n_\mu^2}{2!} - \frac{n_\mu^3}{3!} + \ldots) \hat{n}_\mu] \end{split}$$

(since $\hat{n} = \hat{n}^k$ for all $k \in \mathbb{N}$ for fermions since n = 0, 1)

$$= \prod_{\mu} [\mathbb{1} + (e^{-n_{\mu}} - 1)c_{\mu}^{\dagger}c_{\mu}] \qquad [$$

Let

$$|\phi\rangle = \sum_{i} a_{i} c_{j}^{\dagger} |0\rangle \tag{3.30}$$

be an arbitrary many-particle state.

Now we use the previous identity to prove that applying the operator of equation (3.29) to $|\phi\rangle$ we obtain

$$e^{-c_i^{\dagger} N_{ij} c_j} |\phi\rangle = \sum_j a_j' c_j^{\dagger} |0\rangle, \qquad (3.31)$$

with

$$a'_{j} = \sum_{i} (e^{-B})_{ji} a_{i} \tag{3.32}$$

We start by writing $|\phi\rangle$ in the basis $\{|\mu\rangle\}$ (in which N is diagonal).

$$|\phi\rangle = \sum_{i,\mu} a_i \langle \mu | i \rangle c_{\mu}^{\dagger} | 0 \rangle \tag{3.33}$$

Then, we apply the RHS of equation (3.29) to $|\phi\rangle$ written in this basis.

$$\sum_{\nu} \left[\mathbb{1} + (e^{-n_{\mu}} - 1)c_{\nu}^{\dagger}c_{\nu} \right] c_{\mu}^{\dagger} |0\rangle$$

$$= \left[\mathbb{1} + (e^{-n_{\mu}} - 1)c_{\mu}^{\dagger}c_{\mu} \right] c_{\mu}^{\dagger} |0\rangle$$

$$= c_{\mu}^{\dagger} |0\rangle + (e^{-n_{\mu}-1} - 1)c_{\mu}^{\dagger} |0\rangle$$

$$= c_{\mu}^{\dagger} e^{-n_{\mu}} |0\rangle$$
(3.34)

$$\sum_{\nu} \left[\mathbb{1} + (e^{-n_{\mu}} - 1)c_{\nu}^{\dagger}c_{\nu} \right] |\phi\rangle$$

$$= \sum_{i\mu} \langle \mu | i \rangle a_{i}e^{-n_{\mu}}c_{\mu}^{\dagger}$$

$$= \sum_{j\mu i} \langle j | \mu \rangle e^{-n_{\mu}} \langle \mu | i \rangle a_{i} | j \rangle$$

$$= \sum_{ji} \underbrace{\sum_{\mu\nu} \langle j | \mu \rangle e^{-N_{\mu\nu}} \langle \nu | i \rangle}_{(e^{-N})_{ji}} a_{i} | j \rangle$$

$$= \sum_{j} a'_{j}c_{j}^{\dagger} |0\rangle$$
(3.35)

Similarly, by repeating the procedure performing a change of basis to the eigenbasis of M, we obtain the more general relation

$$e^{-c_i^{\dagger} M_{ij} c_j} e^{-c_i^{\dagger} N_{ij} c_j} |\phi\rangle = \sum_j a_j'' c_j^{\dagger} |0\rangle$$

$$a_j'' = \sum_i (e^{-\mathbf{M}} e^{-\mathbf{N}})_{ji} a_i$$
(3.36)

The amplitude of a propagated state is given by multiplying the initial amplitude by the matrix $e^{-M}e^{-N}$, whichever the basis we choose. Then, since equation (3.36) holds in particular for the choice of the eigenbasis of $e^{-M}e^{-N}$ as our basis of single-particle states, if we start with an eigenstate

$$|\phi\rangle = c_{\nu}^{\dagger} |0\rangle \,, \tag{3.37}$$

then the amplitude of the propagated state will be given by

$$(e^{-M}e^{-N})_{\nu\nu} = e^{-\rho_{\nu}},\tag{3.38}$$

the same as we would obtain from equation (3.24). Clearly, if we start with a state that is an arbitrary combination of states of the eigenbasis, we would obtain the identity (3.24).

The identity was proven for a single-particle state. Does it generalize to more than one particle? As we did before, we start with propagation by a single factor e^{-N} . Take a two-particle state

$$|\phi\rangle = c^{\dagger}_{\mu_1}c^{\dagger}_{\mu_2}\,|0\rangle \tag{3.39}$$

Now propagate it with N, i.e.

$$e^{-c_i^{\dagger} N_{ij} c_j} |\phi\rangle = \prod_{\mu} \left[1 + (e^{-n_{\mu}} - 1) c_{\mu}^{\dagger} c_{\mu} \right] c_{\mu_1}^{\dagger} c_{\mu_2}^{\dagger} |0\rangle$$

$$= e^{-n_{\mu_1}} e^{-n_{\mu_2}} c_{\mu_1}^{\dagger} c_{\mu_2}^{\dagger} |0\rangle,$$
(3.40)

where we simply note that by similar reasoning to the previous case, we would in equation (3.34) keep two terms corresponding to $\mu_1 \neq \mu_2$. If $\mu_1 = \mu_2$, then both sides are equal to zero due to Pauli's

exclusion principle and the equality holds trivially. This reasoning clearly generalizes to an arbitrary superposition of many-particle states. Moreover, we proved the result for a product of two factors $e^{-M}e^{-N}$, but it is also easy to see that by successive changes of basis, we could extend our result to an arbitrary number of factors.

To complete our proof of the identity (3.23) that is so crucial in formulating AFQMC, we use the auxiliar identity we just proved (3.24).

$$\operatorname{Tr}\left[e^{-\sum_{\nu}c_{\nu}^{\dagger}\rho_{\nu}c_{\nu}}\right] = \operatorname{Tr}\left[\prod_{\nu}e^{-c_{\nu}^{\dagger}\rho_{\nu}c_{\nu}}\right] \text{ since } [\hat{n}_{\mu},\hat{n}_{\nu}] = 0$$

$$= \prod_{\nu}(1 + e^{-\rho_{\nu}}) = \operatorname{det}[\boldsymbol{I} + e^{-\boldsymbol{M}}e^{-\boldsymbol{N}}], \qquad \Box$$
(3.41)

where the last equality stems from the fact that the determinant of a diagonal matrix is just the product of the eigenvalues.

When applied to our problem, equation (3.22) essentially makes the computation of the trace possible! Note that if we were to compute it naïvely, we would soon run out of computer memory. The dimension of the Hilbert space of the Hubbard model is exponential in N (actually 4^N), where N is the number of lattice sites. The determinant is calculated for a matrix whose size is polynomial in N.

Equation (3.22) allows us to write the partition function (3.18) in computable form

$$\begin{split} Z_{\boldsymbol{h}} &= \mathrm{Tr}_{\boldsymbol{h}} \bigg[(c_U)^{NL} \mathrm{det}[\boldsymbol{M}_{\uparrow}(\boldsymbol{h})] \mathrm{det}[\boldsymbol{M}_{\downarrow}(\boldsymbol{h})] \bigg] \\ &\equiv \mathrm{Tr}_{\boldsymbol{h}} \bigg[\tilde{\rho}_{\mathrm{eff}}(\boldsymbol{h}) \bigg], \end{split} \tag{3.42}$$

where the fermion matrices M_{σ} are defined in terms of the B-matrices for a given spin σ and a given HS-field h:

$$M_{\sigma}(h) = I + B_{L,\sigma}(h_L)B_{L-1,\sigma}(h_{L-1})...B_{1\sigma}(h_1)$$
 (3.43)

Equation (3.42) defines an effective density matrix (now a function!) $\tilde{\rho}_{\text{eff}}(h)$ in the HS-field space. The computable approximation of the distribution operator \mathcal{P} corresponding to this partition func-

$$P(\mathbf{h}) = \frac{A}{Z_h} \det[\mathbf{M}_{\uparrow}(\mathbf{h})] \det[\mathbf{M}_{\downarrow}(\mathbf{h})], \tag{3.44}$$

where $A = (c_U)^{NL}$ is a normalization constant. This is now a distribution function over configurations h since the problem is classical!

For the particular case of no interactions U=0, we have that $\nu=0$, and $M_{\sigma}(h)$ are constant matrices, independent of the HS-field. The Trotter-Suzuki approximation then becomes exact and the Hubbard Hamiltonian may be simulated exactly after evaluating $M_{\sigma}(h)$ a single time. No updates are required.

As a final remark, note that we managed to map a quantum problem to a classical problem in higher dimension. The degrees of freedom of the quantum problem correspond to the i indices of the

tion is

c-operators. In our formulation, an additional imaginary time slice index *l* was introduced, leading to a mapping that is not specific to the Hubbard model, but that actually applies very generally for any quantum system.

3.1.2.B Monte Carlo sampling of the HS-field

A – The computational problem is now that of sampling configurations of the h field drawn from the distribution P(h) using *Classical* Monte Carlo. The size of the state space has been (hopefully) reduced to 2^{NL} (assuming that L < N).

It remains to choose a dynamics and a sampling scheme. The simplest strategy to change from a configuration h to a new one h' is single spin-flip dynamics. We choose a random point (l,i), and we flip the spin at that "site"

$$h'_{l,i} = -h_{l,i}, (3.45)$$

keeping all others unchanged.

The most common scheme to ensure that the distribution of the accepted sample is P(h) is the Metropolis-Hastings algorithm.

After the warm-up steps, i.e. after we ensure that we are correctly sampling from the required distribution, we may perform measurements, waiting for some (Monte Carlo) time before each of them to ensure that the correlations within the sample are negligible. Let the total number of Monte Carlo steps (warm-up W + measurement M) be S=W+M. The idea is that we run the algorithm for W steps, before starting the measurements. Then we measure the state of the system every 2τ steps, where τ is the correlation time, i.e. the time it takes for some representative correlation function to drop to e^{-1} its original value.

The Metropolis acceptance/rejection scheme leads to a rank-one update of the matrices $M_{\sigma}(h)$, which affords an efficient evaluation of the acceptance ratio $a_{l,i}$ [?].

Consider two matrices A_1 , A_2 written in the form

$$A_{1,2} = I + FV_{1,2}, (3.46)$$

where F is some matrix. $V_{1,2}$ are diagonal and non-singular and differ only in the (1,1) entry, so that

$$V_1^{-1}V_2 = I + \alpha_1 e_1 e_1^T, \tag{3.47}$$

where e_1 is a vector corresponding to the first column of the identity matrix I, and

$$\alpha_1 = \frac{V_2(1,1)}{V_1(1,1)} - 1$$

Then, A_2 is clearly a rank-one update of A_1 .

Algorithm 3.1 Auxiliary Field Quantum Monte Carlo

```
Initialize HS field h
       Initialize hoppings K (h_{l,i}) = (\pm 1)_{l=1,i=1}^{L,N}
       (l,i) \leftarrow (1,1)
 5: for step = 1 to S do
           Propose new configuration by flipping a spin
            h'_{l,i} = -h_{l,i}
            Compute the acceptance ratio a_{l,i}
            \frac{\det[\boldsymbol{M}_{\uparrow}(\boldsymbol{h}')]\det[\boldsymbol{M}_{\downarrow}(\boldsymbol{h}')]}{\det[\boldsymbol{M}_{\uparrow}(\boldsymbol{h})]\det[\boldsymbol{M}_{\downarrow}(\boldsymbol{h})]}
            Metropolis step
            Draw random number r \in [0, 1]
10:
           if r \leq \min(1, a_{l,i}) then
                h = h'
            else
                h = h
            end if
15:
           Next site
           \quad \text{if } i < N \text{ then} \\
                l = l , i = i + 1
            else
                \quad \text{if } l < L \text{ then} \\
                    \boldsymbol{l}=\boldsymbol{l}+\boldsymbol{1} , i=1
20:
                end if
                if l = L then
                   l=1 , i=1
                end if
25:
            end if
       end for
```

$$egin{aligned} m{A}_2 &= m{I} + m{F}m{V}_1 + m{F}m{V}_1 (m{V}_1^{-1}m{V}_2 - m{I}) \ &= m{A}_1 + lpha_1 (m{A}_1 - m{I}) m{e}_1 m{e}_1^T \ &= m{A}_1 [m{I} + lpha_1 (m{I} - m{A}_1^{-1}) m{e}_1 m{e}_1^T] \end{aligned}$$

Using the identity $det[I + xy^T] = 1 + y^Tx$ for any two column vectors, we may write the ratio of the determinants of matrices A_1 and A_2 as

$$r_1 = \frac{\det[A_2]}{\det[A_1]} = 1 + \alpha_1 (1 - e_1^T A_1^{-1} e_1), \tag{3.48}$$

which reduces the computation of the ratio r_1 to computing the (1,1) entry of \mathbf{A}^{-1} .

Now we generalize this idea for a sequence of matrices $A_1, A_2, ..., A_i, ..., A_n$ generated by successive rank-one updates: $A_{i+1} = I + FV_{i+1}, i = 1, 2, ..., n-1$, with

$$V_i^{-1}V_{i+1} = I + \alpha_i e_i e_i^T \quad \alpha_i = \frac{V_{i+1}(1,1)}{V_i(1,1)} - 1$$
 (3.49)

The Sherman-Morrison-Woodbury formula gives an expression for the inverse of A_2 as a rank-one update of A_1^{-1} .

$$\mathbf{A}_{2}^{-1} = \left[\mathbf{I} - \frac{\alpha_{1}}{r_{1}} (\mathbf{I} - \mathbf{A}_{1}^{-1}) \mathbf{e}_{1} \mathbf{e}_{1}^{T} \right] \mathbf{A}_{1}^{T}$$

$$= \mathbf{A}_{1}^{-1} - \frac{\alpha_{1}}{r_{1}} \mathbf{u}_{1} \mathbf{w}_{1}^{T},$$
(3.50)

where

$$u_1 = (I - A_1^{-1})e_1$$
 $w_1 = (A_1^{-1})^T e_1$

Using equation (3.48), we find the updates

$$r_i = \frac{\det[\boldsymbol{M}_{i+1}]}{\det[\boldsymbol{M}_i]} = 1 + \alpha_i (1 - \boldsymbol{e}_i^T \boldsymbol{A}_i^{-1} \boldsymbol{e}_i), \text{ and}$$

$$\boldsymbol{M}_{i+1}^{-1} = \boldsymbol{M}_i^{-1} - \frac{\alpha_i}{r_i} \boldsymbol{u}_i \boldsymbol{w}_i^T,$$
(3.51)

where $\boldsymbol{u}_i = (\boldsymbol{I} - \boldsymbol{A}_i^{-1})\boldsymbol{e}_i$ and $\boldsymbol{w}_i = (\boldsymbol{A}_i^{-1})^T\boldsymbol{e}_i$.

It is possible to generalize this procedure to compute the inverse of M_k as a rank-(k-1) update of A_1^{-1} :

$$M_k^{-1} = M_1^{-1} - U_{k-1} D_k W_{k-1}^T, (3.52)$$

where

$$U_k = [u_1, u_2, ..., u_{k-1}]$$
 and $W = [w_1, w_2, ..., w_{k-1}],$ (3.53)

and $D_k = \operatorname{diag}(\alpha_1/r_1, \alpha_2/r_2, ..., \alpha_{k-1}/r_{k-1}).$

3.1.2.C Making measurements

A- In QMC simulations, physical observables are extracted by measuring them directly over the course of the sampling of the configuration space. The single-particle (equal time) Green's Function is useful to obtain quantities such as density and kinetic energy. It turns out that it is simply the inverse of the M-matrix that we already compute to obtain the acceptance ratio at each step.

$$G_{ij}^{\sigma} = \left\langle c_{i\sigma} c_{j\sigma}^{\dagger} \right\rangle_{\mathbf{h}}$$

$$= \left(M_{\sigma}^{-1}(\mathbf{h}) \right)_{ij}$$

$$= \left([\mathbf{I} + \mathbf{B}_{L,\sigma}(h_L) \mathbf{B}_{L-1,\sigma}(h_{L-1}) ... \mathbf{B}_{1,\sigma}(h_1)]^{-1} \right)_{ij}$$
(3.54)

The equal time Green's function is a fermion average for a given HS-field configuration [?]. The corresponding thermal average is given by

$$\left\langle c_{i}c_{j}^{\dagger} \right\rangle = \frac{1}{Z} \text{Tr} \left[e^{-\beta \mathcal{H}} c_{i} c_{j}^{\dagger} \right]$$

$$= \frac{1}{Z} \text{Tr}_{h} \text{Tr} \left[c_{i\sigma} c_{j\sigma}^{\dagger} \prod_{l=1}^{L} B_{l,\uparrow}(\boldsymbol{h}_{l}) B_{l,\downarrow}(\boldsymbol{h}_{l}) \right],$$

$$(3.55)$$

The density matrix $e^{-\beta \mathcal{H}}$ may be written as a trace over HS-field configurations of a product of L factors corresponding to each imaginary time slice. Recall equation (3.18): the partition function Z is just the trace over the Hilbert space of the aforementioned density matrix. Equivalently, it may

be thought of as a trace over HS-field configurations of the effective density matrix $\hat{\rho}_{\text{eff}}(h)$ defined in equation (3.42).

The Green's function is defined for fixed h. Omitting the spin index σ , without loss of generality, we obtain

$$G_{ij} \equiv \left\langle c_i c_j^{\dagger} \right\rangle_{\mathbf{h}} = \frac{\text{Tr}[\mathbf{B}_L(h_l)\mathbf{B}_{L-1}(h_{l-1})...\mathbf{B}_1(h_1)c_i c_j^{\dagger}]}{\tilde{\rho}_{\text{eff}}}$$

$$= \frac{\text{Tr}[\mathbf{B}_L \mathbf{B}_{L-1}...\mathbf{B}_1 c_i c_j^{\dagger}]}{\text{Tr}[\mathbf{B}_L \mathbf{B}_{L-1}...\mathbf{B}_1]}$$
(3.56)

The trace is evaluated by changing to a basis $\{|\alpha\rangle\}$, where c_i is diagonal and then repeating the procedure for c_j^{\dagger} , now changing again to a basis $\{|\beta\rangle\}$, where c_j^{\dagger} is diagonal. Using equation (3.28), we obtain

$$\left\langle c_{i}c_{j}^{\dagger}\right\rangle_{\mathbf{h}} = \frac{\sum_{\alpha,\beta}\left\langle i|\alpha\right\rangle\left\langle\beta|j\right\rangle\operatorname{Tr}[c_{\alpha}c_{\beta}^{\dagger}\boldsymbol{B}_{L}\boldsymbol{B}_{L-1}...\boldsymbol{B}_{1}]}{\operatorname{Tr}[\boldsymbol{B}_{L}\boldsymbol{B}_{L-1}...\boldsymbol{B}_{1}]}$$
(3.57)

After taking the trace, (on the diagonal basis) the only nonzero contribution will be for $\alpha=\beta$. When $c_{\alpha}c_{\beta}^{\dagger}$ acts on the bra to its left, only that term survives in the sum since c_{α} is a diagonal operator in the basis $\{|\alpha\rangle\}$. On the other hand, the second equality in equation (3.41) gives the contribution to the trace of the exponential of $c_{\alpha}^{\dagger}c_{\alpha}$ appearing in the B-matrices.

$$Tr[B_L B_{L-1}...B_1] = \prod_{\nu} (1 + e^{-\rho_{\nu}}), \tag{3.58}$$

where $\{|\nu\rangle\}$ is the basis in which the product of the B's is diagonal.

$$\left\langle c_{i}c_{j}^{\dagger}\right\rangle_{\mathbf{h}} = \sum_{\alpha} |\alpha\rangle\langle i| \frac{\operatorname{Tr}[c_{\alpha}c_{\alpha}^{\dagger}B_{L}B_{L-1}...B_{1}]}{\operatorname{Tr}[B_{L}B_{L-1}...B_{1}]} |j\rangle\langle\alpha|$$

$$= \sum_{\alpha} |\alpha\rangle\langle i| \frac{\operatorname{Tr}[(1-c_{\alpha}^{\dagger}c_{\alpha})B_{L}B_{L-1}...B_{1}]}{\operatorname{Tr}[B_{L}B_{L-1}...B_{1}]} |j\rangle\langle\alpha|$$

$$= \sum_{\alpha} |\alpha\rangle\langle i| 1 - \frac{\operatorname{Tr}[c_{\alpha}^{\dagger}c_{\alpha}e^{-\Delta\tau\hat{h}}]}{\operatorname{Tr}[e^{-\Delta\tau\hat{h}}]} |j\rangle\langle\alpha|$$

$$= \sum_{\alpha} |\alpha\rangle\langle i| 1 - \frac{1}{1+e^{\Delta\tau\varepsilon_{\alpha}}} |j\rangle\langle\alpha|$$

$$= \sum_{\alpha} |\alpha\rangle\langle i| \frac{e^{\Delta\tau\varepsilon_{\alpha}}}{1+e^{\Delta\tau\varepsilon_{\alpha}}} |j\rangle\langle\alpha|$$

$$= \sum_{\alpha} |\alpha\rangle\langle i| \frac{1}{1+e^{-\Delta\tau\varepsilon_{\alpha}}} |j\rangle\langle\alpha|$$

$$= \sum_{\alpha} |\alpha\rangle\langle i| \frac{1}{1+e^{-\Delta\tau\varepsilon_{\alpha}}} |j\rangle\langle\alpha|$$

$$= \left[\frac{1}{B_{L}B_{L-1}...B_{1}}\right]_{ij},$$

$$(3.59)$$

where in the fourth equality we used an analogy with the Fermi function defined as

$$f_{\alpha} = \frac{\text{Tr}[e^{-\beta \mathcal{H}} \hat{n}_{\alpha}]}{\text{Tr}[e^{-\beta \mathcal{H}}]} = \left(1 + e^{\beta \varepsilon_{\alpha}}\right)^{-1} \tag{3.60}$$

for $\mu=0$ and with $\beta\mapsto\Delta\tau$. The product of \boldsymbol{B} -matrices was written as the exponential $e^{-\Delta\tau\hat{h}}$, which can be done because we have shown before that it is possible to diagonalize the product in a basis in which the trace amounts to the simple form of equation (3.58).

An alternative way of arriving to this result is to note that in the expression we obtain in the second equality, only the term $\nu=\alpha$ from equation (3.58) contributes [?], leading to the final result with $\rho_{\alpha}=\Delta \tau \varepsilon_{\alpha}$.

The electron density may be obtained from the Green function

$$\rho_{i\sigma} = \left\langle c_{i\sigma}^{\dagger} c_{i\sigma} \right\rangle = 1 - \left\langle c_{i\sigma} c_{i\sigma}^{\dagger} \right\rangle = 1 - G_{ii}^{\sigma}, \tag{3.61}$$

It is natural to think of averaging this over the lattice, and over the spins. This is justified by the fact that the Hubbard Hamiltonian is translationally invariant. Thus, $\rho_{i\sigma}$ should be independent of the spatial site. This statement is strict when exactly solving the model, but it becomes only approximate, i.e. valid only on average in our simulations. Thus, we take the average

$$\rho = \frac{1}{2N} \sum_{\sigma} \sum_{i=1}^{N} \rho_{i\sigma} \tag{3.62}$$

in an attempt to reduce statistical errors.

One must pay attention to the symmetry of the model at hand, since a similar model for a disordered system including randomness would not be translationally invariant anymore. Moreover, it is implicit that $\rho_{i\sigma}$ is already averaged over the HS-field configurations that were sampled through the simulation.

The average kinetic energy is similarly obtained.

$$\langle \mathcal{H}_K \rangle = -t \sum_{\langle i,j \rangle, \sigma} \left\langle (c_{i\sigma}^{\dagger} c_{j\sigma} + c_{j\sigma}^{\dagger} c_{i\sigma}) \right\rangle$$

$$= t \sum_{\langle i,j \rangle, \sigma} (G_{ij}^{\sigma} + G_{ji}^{\sigma}), \tag{3.63}$$

where the minus sign is due to the switching of the order of the operators bringing the c^{\dagger} to the right.

3.1.2.D Correlation functions

A – One of the most important goals of QMC simulations is to inspect the system for order of various types, and to find associated phase transitions. This is done by computing correlation functions C(j), measuring how correlated two sites separated by a distance j are.

$$C(j) = \left\langle \mathcal{O}_{i+j} \mathcal{O}_i^{\dagger} \right\rangle - \left\langle \mathcal{O}_{i+j} \right\rangle \left\langle \mathcal{O}_i^{\dagger} \right\rangle, \tag{3.64}$$

where \mathcal{O} is an operator corresponding to the order parameter of the phase transition. For example, we might be looking for magnetic order, in which case the relevant operators are $\mathcal{O}_i = n_{i\uparrow} - n_{i\downarrow}$, $\mathcal{O}_i^{\dagger} = n_{i\uparrow} - n_{i\downarrow}$, or superconductivity, where we would like to measure correlations in fermion pair formation: $\mathcal{O}_i = c_{i\downarrow}c_{i\uparrow}$, $\mathcal{O}_i^{\dagger} = c_{i\uparrow}^{\dagger}c_{i\downarrow}^{\dagger}$.

In general, we expect a high temperature disordered phase, for which correlations decay exponentially $C(j) \propto e^{-j/\xi}$, where ξ is a characteristic length called the correlation length. At some point, there can be a transition to a low temperature phase, where $C(j) \propto m^2$, where m is the order parameter for the transition. Right at the transition, that is at $T=T_c$, there might be singular behavior. In continuous phase transitions, the correlation length diverges $\xi \propto (T-T_c)^{-\nu}$, and the correlations decay slower (in fact algebraically): $C(j) \propto j^{-\eta}$, in an intermediate behavior between exponential decay and a constant. The *critical* exponents ν , and η are characteristic of the transition, or more accurately, of the universality class it belongs to.

The behavior of all these quantities on finite lattices does not precisely correspond to the infinite system behavior. The tails of the functions, i.e. the $j \to \infty$ limit is not well captured. Finite-size scaling is a method to improve on these predictions.

To evaluate correlation functions we use Wick's theorem. Expectations of more than two fermion creation and annihilation operators reduce to products of expectations of pairs of creation and annihilation operators. For example, for pair order:

$$\langle C(j) \rangle = \langle c_{i+j,\downarrow} c_{i+j,\uparrow} c_{i,\uparrow}^{\dagger} c_{i,\downarrow}^{\dagger} \rangle \tag{3.65}$$

How would one measure a correlation function experimentally? Fortunately, there is a quantity that is easy to measure called structure factor, which is just the Fourier transform of the correlation function

$$S(q) = \sum_{j} e^{iqj} C(j) \tag{3.66}$$

The accuracy of QMC simulations can be evaluated by comparing the results for correlation functions with the corresponding structure factors, which can be measured experimentally.



Figure 3.1: Dummy Figure Caption.

Remember you can change the reference style. Another dummy citation [?].

3.2 Section B

3.2.1 Subsection A

The model described can also be represented as

$$\dot{\mathbf{x}}(t) = \mathbf{T}\mathbf{z}(y), \ \mathbf{y}(0) = \mathbf{y}_0, \ z \ge 0$$
 (3.67)

where

$$\mathbf{A} = \begin{bmatrix} -(a_{12} + a_{10}) & a_{21} \\ a_{12} & -(a_{21} + a_{20}) \end{bmatrix}, \ \mathbf{x} = \begin{bmatrix} x_1 \\ x_2 \end{bmatrix}$$
 (3.68)

3.2.2 Subsection B

Table 3.1: Dummy Table.

Vendor Name	Short Name	Commercial Name	Manufacturer	
Text in Multiple Row	ABC	ABC [®]	ABC SA	
	DEF	DEF®	DEF SA	
	GHF	GHF [®]	GHF SA	
Text in Single Row	IJK	IJK [®]	IJK SA	
Frescos SA	LMN	LMN®	LMN SA	
Carros Lda.	Text in Multiple Column			

Conclusions and Future Work

Conclusions Chapter

Title of AppendixA