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First Year Report

The Origin and Significance of the Phase of out of plane Exchange Coupling

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List of Abbreviations

2DBZ	two dimensional Brillouin zone
AF	anti-ferromagnetic
bcc	body-centre cubic
CIP	current in-plane
CPP	current perpendicular to plane
CSP	Cunningham Special Points
fcc	face-centre cubic
FFT	fast Fourier transform
FM	ferromagnetic
GMR	giant magneto-resistance
HDD	hard disk drive
IEC	interlayer exchange coupling
LDOS	local density of states
MRAM	magnetoresistive random access memory
MTGF	Möbius-transformed Green's function
MTJ	magnetic tunnel junctions
PM	polarising magnet
SC	simple cubic
SGF	surface Green's function
SM	switching magnet
SPA	stationary phase approximation
STT	spin transfer torque
TMR	tunnelling magneto-resistance

1 Author's developments since registration

The following programming languages have been studied, and listed in order of fluency. Libraries studied are given in sub-entries. The references given indicate the documentation studied.

- C++98 [\[1,2\]](#), C++11 [\[3\]](#), C++14
 - Eigen [\[4\]](#)
 - GNU Scientific Library [\[5\]](#)
 - OpenMP [\[6\]](#)
 - Blaze [\[7\]](#)
 - Boost
 - SymbolicC++ [\[8\]](#)
 - CBLAS
 - CLAPACK
- Julia [\[9\]](#)
 - SymPy
- Python2.7 [\[10\]](#), Python3.4
 - SymPy [\[11\]](#)
 - SciPy [\[12\]](#)
 - NumPy [\[13\]](#)
- Fortran95
- C
 - GNU Scientific Library
 - CBLAS
 - CLAPACK
- Maple 2015 programming blocks
- Scilab [\[14\]](#)
- Octave [\[15\]](#)
- Haskell [\[16\]](#)
- Perl
- Go

The following software packages have been studied with a brief description of their purposes

- Maple 2015 - *computer algebra package*
- VIM [\[17\]](#) - *text editor with powerful key-bindings*
 - Latex-Suite - *L^AT_EX package with additional key-bindings specific to L^AT_EX*
- TeXMaker - *a L^AT_EX editor*
- QTIplots - *provides graphical output from raw data*
- L^AT_EX - *a powerful typesetting language*
 - TikZ [\[18\]](#) - *provides graphics through a coding interface; all figures in this document are TikZ images*
 - pgfplots [\[19\]](#) - *built from TikZ, a package that generates graphs within L^AT_EX; all graphs in this document are built using pgfplots*
 - BibT_EX - *a bibliographical package*
 - glossaries - *generates glossaries and lists of acronyms if required, by running a Perl script*

Further skills developed include learning to touch-type. Miscellaneous tools used include the LLVM Clang compiler for C++, the GNU Debugger, Google gperftools and valgrind. Git [\[20\]](#) is used for version control, remote repositories are stored on the github server under the user adurie. The following programs have been written with a brief description of their purposes:

- C++
 - cunningham.h - *header file that performs two dimensional Brillouin zone (2DBZ) integration. Further details are given in the Appendix*
 - cunningham_spawn.h - *header file that performs adaptive 2DBZ integration. Further details are given in the Appendix*
 - interp.cpp - *takes raw data as input and interpolates results*
 - analytic_conductance_trilayer_landauer.cpp - *calculates the conductance using the Landauer formalism using analytical methods*
 - numerical_conductance_trilayer_landauer.cpp - *calculates the conductance using the Landauer formalism using direct calculations*
 - LDOS_SGF.cpp - *calculates the local density of states (LDOS) of a semi-infinite crystal at user-specified atomic planes from the surface, using Equation (55)*
 - trilayer_conductance_kubo.cpp - *calculates the charge current using the Kubo formula [\[21\]](#) in the spinless one-band simple cubic (SC) model*
 - spincurrent_4-layer_<j_n-1>2_keldysh.cpp - *calculates the transport term of the spincurrent (Equation (89)) under the Keldysh Formalism [\[22\]](#) in one-band SC. The integration over energy has been bypassed by observing the first term of the Taylor expansion of the trailing Fermi terms is a delta-function. At infinitesimal bias, here $V_b = 1\text{eV}$, this first term is the only term non-negligible*
 - spinenergy.cpp - *calculates the equilibrium term, Equation (88) of the spincurrent under the Keldysh formalism in the one-band SC system. Difficulties with convergence of the energy integral, as yet unresolved*
 - EC_integration_GSL.cpp - *calculates the interlayer exchange coupling (IEC) in one-band SC by calculating Equation (94) directly. Convergence problems not yet resolved*
 - matsubara_method_EC_trilayer.cpp - *calculates the IEC in one-band SC using Equation (98)*
 - EC_SPA_final_efficient.cpp - *calculates the IEC in one-band SC using the*

stationary phase approximation (SPA) (Equation (123)) where the Fourier components are obtained directly. Utilises OpenMP for concurrency

- EC_SPA_FFT.cpp - *calculates the IEC in one-band SC using the SPA (Equation (123)) where the Fourier coefficients are obtained through an inverse fast Fourier transform (FFT). Utilises OpenMP for concurrency*
- EC_SPA_mobius.cpp - *calculates the IEC in one-band SC using the SPA (Equation (123)) where the Fourier coefficients are obtained from Equation (132). Utilises OpenMP for concurrency*
- EC_2-band.cpp - *calculates the IEC in a two-band trilayer using Equation (98)*
- Fortran
 - conductor_approx_with_r.f90 - *calculates the charge current for a trilayer with a conducting spacer using the Landauer formalism*
 - insulator_approx_with_r^2.f90 - *calculates the charge current for a trilayer with an insulator spacer using the Landauer formalism*
 - analytic_trilayer_with_conductor.f90 - *calculates the charge current for a trilayer with a conducting spacer using analytic derivations using the Landauer formalism*

In addition to the equations derived in this document, the following known results have also been derived by the author:

- Transport spincurrent Keldysh term without the energy integral, derived from Equation (89)
- Partial conductance of a trilayer with an insulating spacer, Equation (3) of [23]
- Analytic form of total conductance in a trilayer with an insulating spacer using the Landauer formalism, Equation (6) of [23]
- Analytic form of total conductance in a trilayer with a conducting spacer using the Landauer formalism by considering Γ -point and cutoff-point terms using higher order stationary phase approximations and cutoff point approximations and further Fourier terms. These were written up in an unpublished document [24]
- Partial conductance of a 4-layer system in the Landauer formalism in a conductor
- Creation of a Maple document that calculates the partial conductance in the Landauer formalism in a conductor for an arbitrary number of layers using a programming block
- Further tight binding
 - calculation of Fermi surfaces of SC, face-centre cubic (fcc) and body-centre cubic (bcc)

- SC 2-band model
- fcc in a two atom basis in a pseudo one-dimensional lattice

A great deal of the programs written and formulae derived were done in conjunction with Maple worksheets.

Abstract

The system explored is a pseudo one-dimensional magnetic multilayer. The Green's function formalism is used here, as this offers a means of providing complete and accurate results comparable to experiment, if correctly modelled. This review begins with a brief history of the subject area, moving on to a description of the necessary tools needed, such as the method of Tight Binding and derivations of closed form Green's functions. These methods are then applied to the IEC in its full form simplified by the method of Matsubara, and followed by the derivation of an analytical method. The origin of the phase of IEC is then discussed and derived to be followed by a conclusion and a plan for future work. The paper concludes with numerical integration routines described in the Appendix.

2 Introduction

2.1 History

For the past three decades, there has been significant interest in the field of spintronics. This interest began with the discovery that the resistance of a multilayer comprising of alternating sets of thin magnetic and equally thin non-magnetic metallic layers was much lower when the magnetic moments of the layers were parallel compared to the anti-parallel case [25–27]. This is explained by the fact that if a device is constructed on the nanoscopic scale, such that the thickness is smaller than the spin diffusion length, spin can flow through this device with some memory of its orientation. A magnet by definition polarises spin, so that spincurrent (spin-polarised current) can be modelled as if flowing through two independent wires. The up and down spin carriers are scattered at different rates at layer interfaces and the total resistance of the structure depends on the magnetic configuration of all magnetic components, which in turn can be modified by an applied magnetic field. Thus the resistance can be manipulated by this field. This is giant magneto-resistance (GMR). The so-called *optimistic magneto-resistance ratio* is defined by

$$\frac{\Delta R}{R} = \frac{R^{\uparrow\downarrow} - R^{\uparrow\uparrow}}{R^{\uparrow\uparrow}} \quad (1)$$

Where $R^{\uparrow\uparrow}$, $R^{\uparrow\downarrow}$ are the resistances of the saturated ferromagnetic (FM) configuration and anti-ferromagnetic (AF) configuration respectively. A second significant finding in the field was the discovery by Parkin *et. al.* [28] of IEC which describes the fact that the relative orientation of magnetic moments of coupled magnetic layers separated by a nonmagnetic spacer, is dependent on the spacer layer thickness and in fact oscillates between FM and AF regimes as a function of spacer layer thickness. The early GMR

devices were built with current in-plane (CIP) and were quickly developed into commercial devices, with IBM in 1997 producing hard disk drive (HDD) read heads that utilised this effect to allow a new upper bound on the number of bits read/writable per unit area. In fact, the storage density increased from 0.1 to 100Gbit/in² between 1991 and 2003 [23].

In 2001, two groups simultaneously described the theory of a new system [29,30] which is designed with current perpendicular to plane (CPP). Their original system was a trilayer with ferromagnetic leads made from Fe, and an insulating spacer material made from crystalline MgO. From a theoretical standpoint, these materials are excellent to study since their Brillouin zones are of very similar sizes, meaning growth from one material to another is easy to model since only a change in potentials is required. The magneto-resistance produced by this scheme is known as tunnelling magneto-resistance (TMR) since the electrons must tunnel through the MgO potential barrier, making this an entirely quantum mechanical effect. GMR predicts an optimistic magneto-resistance ratio as high as 50% whereas TMR predicts a ratio as high as 1000%. The success of the TMR system, which was confirmed experimentally in 2004 [31,32], led to a rapid appearance of these TMR read heads in commercial HDDs, making their debut appearance in 2007. As a result, these read heads are now constructed in every modern HDD. The 2004 experiments were unable to replicate the same MR ratio reaching only as high as 200% at room temperature. This is due to the difficulty in engineering 'clean' epitaxial surfaces free from defect, to which the original magnetic tunnel junctions (MTJ) were sensitive. In 2008 experimentalists were able to get TMR ratios as high as 1144% at 5K [33].

It has been well known for some time that altering the magnetic configuration of a magnetic multilayer influences the charge current flowing through it. More recently it has been recognised that conversely, by passing a strong charge current through a magnetic trilayer, the current becomes spin-polarised when passing through the pinned magnet so that a spincurrent flows through the spacer to the opposing ferromagnet (the switching magnet). The spincurrent is absorbed by the switching magnet (SM) exerting a force known as spin transfer torque (STT). If the STT is high enough, the SM undergoes a total reversal of magnetisation in an effect known as current-induced switching of magnetisation. This remains a very active field of research, due to its obvious commercial potential in magnetoresistive random access memory (MRAM) [34], a non-volatile form of random access memory with potential to become a universal memory [35], a form of data storage replacing the separate CPU cache, RAM and HDD storage. MRAM is non-volatile since the SMs retain their configuration in the absence of charge current. It remains a problem that too much heat is generated in these devices during the writing stage, and this is of current research interest.

In the absence of charge current in a magnetic multilayer, there remains a net flow of spincurrent between the two magnetic layers, provided the magnets are out of alignment. This is IEC. Given the phenomenon exists without the expense of a charge current, IEC has been under active research due to possible implications such as MRAM or CPU chips without the same heat generation due to lack of flowing charge. Recently IEC has fallen out of favour, with later research focussing on STT and current-induced

switching of magnetisation. In 1998, Ebels *et. al.* found by experiment induced phase shifts in IEC in magnetic trilayers, by doping the magnetic layers [36]. Due to limited interest at the time, it is likely the potential implications such as switching induced by altering the phase of IEC were not recognised. Recently, one of the pioneers of CPP TMR devices made this observation and revives the mechanism for doing so [37]. This is the main research focus here.

3 Green's Function Formalism

3.1 Method of Tight Binding

The Method of Tight Binding is a convenient framework to use in this context as it models very well the potentials arising in atoms in a regular periodic structure. In an infinite periodic structure, a translation from one crystal site to another leaves the structure as a whole unchanged since each site is identical. A structure of this nature is referred to as Bloch-like [38]. The position vector of each atomic site is given by \mathbf{R} , where $\mathbf{R} = \sum_j^{\chi} n_j \mathbf{a}_j$ can describe the position of any lattice site with integers n_j , primitive translation vectors \mathbf{a}_j and χ is equal to the dimension of \mathbf{R} .

In the tight binding formalism, the energy eigenvalues are given by

$$\varepsilon(\mathbf{k}) = \frac{\langle \mathbf{k} | \mathbf{H} | \mathbf{k} \rangle}{\langle \mathbf{k} | \mathbf{k} \rangle} \quad (2)$$

In order to diagonalise the matrix Hamiltonian, \mathbf{H} , the transformation from real space to momentum space (k -space) is required

$$|\mathbf{k}\rangle = \frac{1}{\sqrt{N}} \sum_{\mathbf{R}=-\infty}^{\infty} e^{i\mathbf{k}\cdot\mathbf{R}} |\mathbf{R}\rangle \quad (3)$$

where the summation is over every lattice point, N is the total number of lattice points in the crystal. The position of any site in k -space is given by $\mathbf{K} = \sum_i^{\chi} k_i \mathbf{b}_i$ with integers k_i and primitive translation vectors \mathbf{b}_i determined such that

$$\mathbf{b}_i \cdot \mathbf{a}_j = 2\pi\delta_{ij} \quad (4)$$

where since

$$\langle \mathbf{R} | \mathbf{R}' \rangle = \delta_{\mathbf{R}\mathbf{R}'} \quad (5)$$

it follows that

$$\langle \mathbf{k} | \mathbf{k}' \rangle = \frac{1}{N} \sum_{\mathbf{R}\mathbf{R}'} e^{-i\mathbf{k}\cdot\mathbf{R}} e^{i\mathbf{k}'\cdot\mathbf{R}'} \langle \mathbf{R} | \mathbf{R}' \rangle \equiv \frac{1}{N} \sum_{\mathbf{R}} e^{-i\mathbf{R}\cdot(\mathbf{k}-\mathbf{k}')} = \delta_{\mathbf{k}\mathbf{k}'} \quad (6)$$

Assuming a small set of nearby atoms overlap with non-negligible results, all elements of the Hamiltonian are zero, except

$$\langle \mathbf{R} | \mathbf{H} | \mathbf{R} \rangle = \mathbf{u} \quad (7)$$

the on-site potential, and

$$\langle \mathbf{R} | \mathbf{H} | \mathbf{R} \pm \mathbf{d} \rangle = \mathbf{t} \quad (8)$$

where \mathbf{d} is the set of position vectors of the non-negligible neighbouring terms. \mathbf{t} is the set of hopping matrices whose strength varies depending on the binding distance to the on-site term. Models involving multiple energy bands or multiple bases per site require that \mathbf{t}, \mathbf{u} are matrices.

Thus

$$\langle \mathbf{R} | \mathbf{H} | \mathbf{k} \rangle = \frac{1}{\sqrt{N}} \sum_{\mathbf{R}} \langle \mathbf{R} | \mathbf{H} | \mathbf{R} \rangle e^{i\mathbf{k} \cdot \mathbf{R}} = \frac{1}{\sqrt{N}} e^{i\mathbf{k} \cdot \mathbf{R}} \left(\mathbf{u} + \sum_{\mathbf{d} \in \mathbf{R}} e^{i\mathbf{k} \cdot \mathbf{d}} \langle \mathbf{R} | \mathbf{H} | \mathbf{R} \pm \mathbf{d} \rangle \right) \quad (9)$$

and

$$\langle \mathbf{k} | \mathbf{H} | \mathbf{k} \rangle = \frac{1}{N} \sum_{\mathbf{R}} e^{i\mathbf{R} \cdot (\mathbf{k} - \mathbf{k})} \left(\mathbf{u} + \sum_{\mathbf{d} \in \mathbf{R}} e^{i\mathbf{k} \cdot \mathbf{d}} \langle \mathbf{0} | \mathbf{H} | \mathbf{d} \rangle \right) = \delta_{\mathbf{k}\mathbf{k}} \left(\mathbf{u} + \sum_{\mathbf{d} \in \mathbf{R}} e^{i\mathbf{k} \cdot \mathbf{d}} \langle \mathbf{0} | \mathbf{H} | \mathbf{d} \rangle \right) \quad (10)$$

where Equation (6) is used. The energy eigenvalues $\varepsilon(\mathbf{k})$ are thus the eigenvalues of

$$\mathbf{u} + \sum_{\mathbf{d} \in \mathbf{R}} e^{i\mathbf{k} \cdot \mathbf{d}} \langle \mathbf{0} | \mathbf{H} | \mathbf{d} \rangle \quad (11)$$

3.1.1 One-band model 1D example

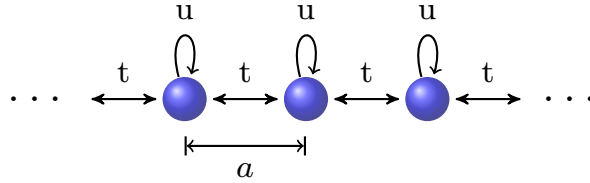


Figure 1: Depiction of an infinite chain of atoms in one dimension separated by a showing the self energy u , and the hopping energy t .

In the one-band model in one dimension, all the vector values become scalar and the position of a given lattice site is $R = na$ and $d = \pm a$. As such, Equation (11) becomes

$$\varepsilon = u + t(e^{ika} + e^{-ika}) = u + 2t \cos ka \quad (12)$$

3.1.2 One-band model bcc example

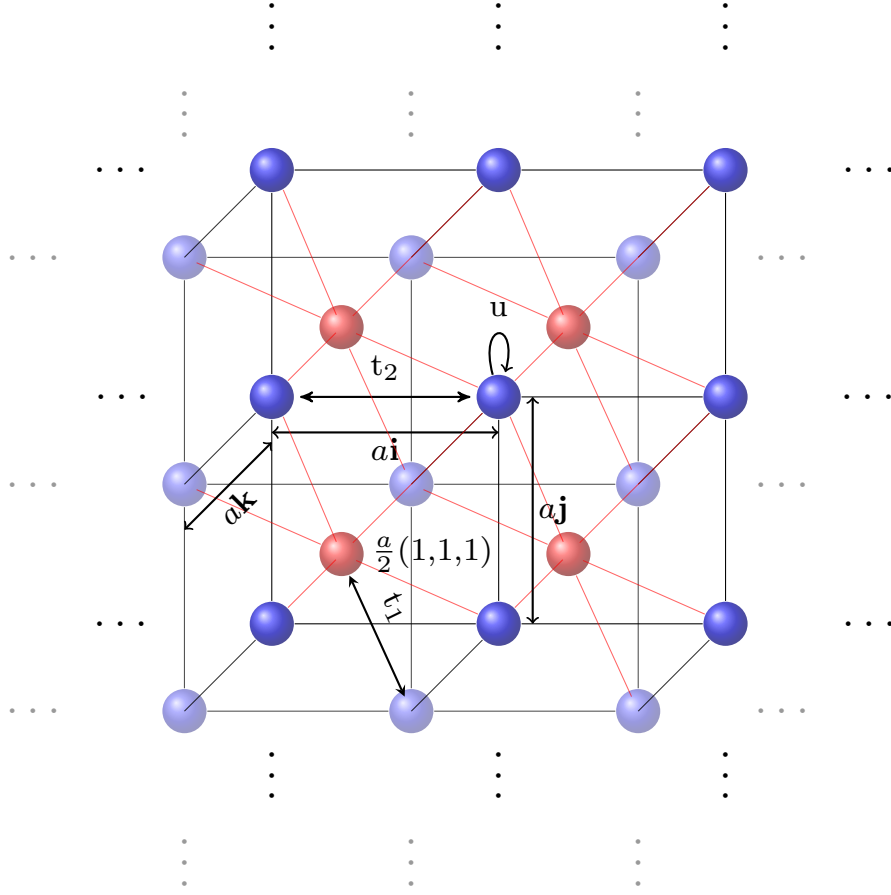


Figure 2: Depiction of an infinite bcc lattice with lattice vector \mathbf{a} . The self energy is u , the nearest neighbour hopping t_1 and the next-nearest neighbour hopping t_2 .

In the bcc structure, the nearest neighbours are the sites in the body centres, however the next nearest neighbours are only 15% further away and so the accuracy is greatly improved when next-nearest neighbour hopping is included [39]. In the one band model, the on-site potential and hoppings are scalar.

Here $\mathbf{R} = n_1\mathbf{a}_1 + n_2\mathbf{a}_2 + n_3\mathbf{a}_3$ with integers n_1, n_2, n_3 and the primitive translation vectors may be chosen as follows;

$$\mathbf{a}_1 = \frac{a}{2}(1, 1, 1), \quad \mathbf{a}_2 = \frac{a}{2}(-1, 1, 1), \quad \mathbf{a}_3 = \frac{a}{2}(1, -1, 1). \quad (13)$$

The position of any site in the reciprocal lattice is $\mathbf{K} = k_1\mathbf{b}_1 + k_2\mathbf{b}_2 + k_3\mathbf{b}_3$ with integers k_1, k_2, k_3 and primitive translation vectors $\mathbf{b}_1, \mathbf{b}_2, \mathbf{b}_3$.

A general procedure for satisfying Equation (4) in 3D is as follows [40]

$$\mathbf{b}_1 = \frac{2\pi(\mathbf{a}_2 \times \mathbf{a}_3)}{\mathbf{a}_1 \cdot (\mathbf{a}_2 \times \mathbf{a}_3)} \quad \mathbf{b}_2 = \frac{2\pi(\mathbf{a}_3 \times \mathbf{a}_1)}{\mathbf{a}_2 \cdot (\mathbf{a}_3 \times \mathbf{a}_1)} \quad \mathbf{b}_3 = \frac{2\pi(\mathbf{a}_1 \times \mathbf{a}_2)}{\mathbf{a}_3 \cdot (\mathbf{a}_1 \times \mathbf{a}_2)} \quad (14)$$

Thus

$$\mathbf{b}_1 = \frac{2\pi}{a}(1, 1, 0) \quad \mathbf{b}_2 = \frac{2\pi}{a}(-1, 0, 1) \quad \mathbf{b}_3 = \frac{2\pi}{a}(0, -1, 1) \quad (15)$$

It is apparent from Equation (15) that the reciprocal lattice forms a fcc structure. It is immediately obvious from studying Equation (10) that the energy eigenvalue in this instance is

$$\varepsilon = u + t_1 \sum_{\mathbf{d} \in \text{nn}} e^{i\mathbf{k} \cdot \mathbf{d}} + t_2 \sum_{\mathbf{d} \in \text{nnn}} e^{i\mathbf{k} \cdot \mathbf{d}} \quad (16)$$

where the left-hand summation is for \mathbf{d} in the set of nearest neighbours (nn), the right-hand summation is for \mathbf{d} in the set of next-nearest neighbours (nnn). t_1, t_2 are the nearest neighbour and next-nearest neighbour hoppings respectively. In the case of the bcc structure described here in the (100) configuration, $\mathbf{d} \in \text{nn}$ is the set of vectors

$$\begin{aligned} \pm \frac{a}{2}(1, 1, 1) &= \pm \mathbf{a}_1, & \pm \frac{a}{2}(-1, 1, 1) &= \pm \mathbf{a}_2, \\ \pm \frac{a}{2}(1, -1, 1) &= \pm \mathbf{a}_3, & \pm \frac{a}{2}(1, 1, -1) &= \pm(\mathbf{a}_1 - \mathbf{a}_2 - \mathbf{a}_3) \end{aligned}$$

and $\mathbf{d} \in \text{nnn}$ is the set of vectors

$$\pm a(1, 0, 0) = \pm(\mathbf{a}_1 - \mathbf{a}_2), \quad \pm a(0, 1, 0) = \pm(\mathbf{a}_1 - \mathbf{a}_3), \quad \pm a(0, 0, 1) = \pm(\mathbf{a}_2 + \mathbf{a}_3)$$

since this structure has eight nearest neighbours and six next-nearest neighbours. Solving Equation (16) and pairing the positive and negative exponents yields

$$\begin{aligned} \varepsilon = u + 2t_1 \left(\cos \frac{a}{2}(k_x + k_y + k_z) + \cos \frac{a}{2}(-k_x + k_y + k_z) + \cos \frac{a}{2}(k_x - k_y + k_z) + \right. \\ \left. \cos \frac{a}{2}(k_x + k_y - k_z) \right) + 2t_2 (\cos k_x a + \cos k_y a + \cos k_z a) \quad (17) \end{aligned}$$

which after some algebra and the use of trigonometric sum and difference identities and a number of cancellations becomes

$$\varepsilon = u + 8t_1 \cos \left(\frac{k_x a}{2} \right) \cos \left(\frac{k_y a}{2} \right) \cos \left(\frac{k_z a}{2} \right) + 2t_2 (\cos k_x a + \cos k_y a + \cos k_z a) \quad (18)$$

with $t_2 < t_1$, since it is part of the definition of tight-binding that the next-nearest neighbours have a weaker interaction than the nearest neighbours.

3.1.3 Two-band model 2D square example

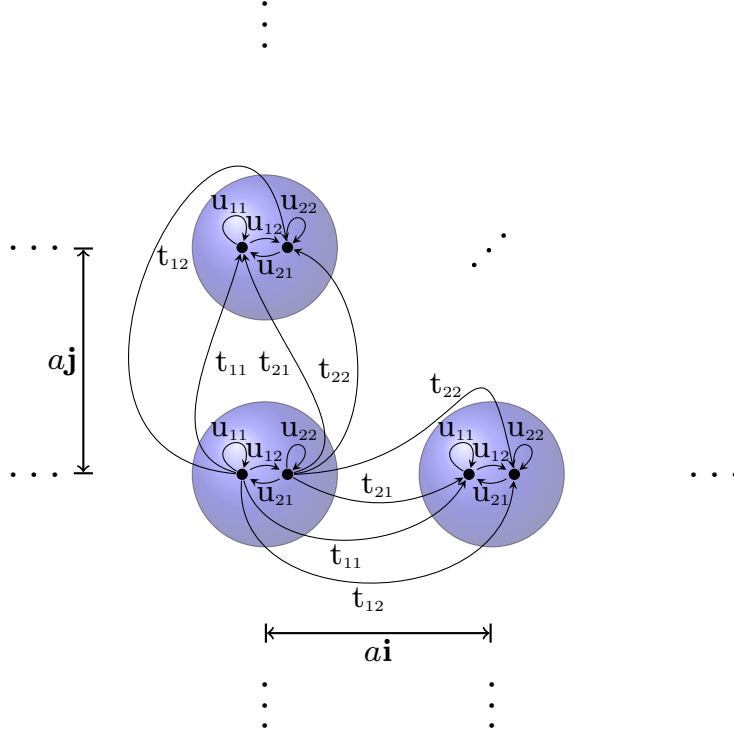


Figure 3: Depiction of an infinite square mesh of atoms in two bands separated by a showing the various energies contained within the matrices in Equation (19). Full details of these energies are given below.

In the two band model, the interaction between electrons at different excited states occupying the two available bands per atomic site must also be considered. This means the potential of an electron associated with the on-site atom is now more complicated. Similarly, hopping to another atomic site is more complicated since there are multiple bands to hop to. Figure 3 indicates the origin of these terms. Nearest neighbour interactions are the only terms considered in this example, in this case Equations (7) & (8) are the 2×2 matrices

$$\mathbf{u} = \begin{pmatrix} u_{11} & u_{12} \\ u_{21} & u_{22} \end{pmatrix}, \quad \mathbf{t} = \begin{pmatrix} t_{11} & t_{12} \\ t_{21} & t_{22} \end{pmatrix} \quad (19)$$

The self-energy is u_{ii} for $i \in [1, 2] \in \mathbb{Z}$ the two bands, and u_{12} (u_{21}) is the energy required to raise (lower) an electron to a more (less) excited state on the same atom. t_{ii} is the energy required to hop to an equivalent excited state on a neighbouring atom, and t_{12} (t_{21}) is the energy required to hop to a more (less) excited state on a neighbouring atom.

In the case of the square lattice, the primitive position vectors are $\mathbf{a}_1 = a(0, 1)$, $\mathbf{a}_2 = a(1, 0)$ and \mathbf{d} is the set of nearest neighbours $\pm\mathbf{a}_1, \pm\mathbf{a}_2$. Thus using Equation (11),

the eigenvalues of the two band system are given by

$$\det \begin{pmatrix} u_{11} + 2t_{11} (\cos(k_x a) + \cos(k_z a)) - \varepsilon(\mathbf{k}) & u_{12} + 2t_{12} (\cos(k_x a) + \cos(k_z a)) \\ u_{21} + 2t_{21} (\cos(k_x a) + \cos(k_z a)) & u_{22} + 2t_{22} (\cos(k_x a) + \cos(k_z a)) - \varepsilon(\mathbf{k}) \end{pmatrix} = 0 \quad (20)$$

thus

$$\varepsilon(\mathbf{k}) = \frac{1}{2}(u_{11} + u_{22}) + (t_{11} + t_{22}) (\cos(k_x a) + \cos(k_z a)) \pm \left[(t_{11}^2 - 2t_{11}t_{22} + 4t_{12}t_{21} + t_{22}^2) (\cos(k_x a) + \cos(k_z a))^2 + \frac{1}{4} (u_{11} - u_{22})^2 + u_{12}u_{21} + ((t_{11} - t_{22})(u_{11} - u_{22}) + 2(t_{12}u_{21} + t_{21}u_{12})) (\cos(k_x a) + \cos(k_z a)) \right]^{\frac{1}{2}} \quad (21)$$

whereupon switching off $u_{12}, u_{21}, t_{12}, t_{21}$ and making $u_{11} = u_{22}, t_{11} = t_{22}$ yields the one-band case, as expected.

3.2 Surface Green's function

The tight binding method is useful for describing infinite crystals, but in the study of multilayer systems, edge effects (surfaces) must be appropriately considered. Surface Green's functions (SGFs) allow this to be possible and are an extension of the tight binding framework. Two methods for determining closed-form SGFs will now be given. The first method, the Kalkstein & Soven method is for historical purposes and has little practical significance as it fails to determine SGFs for more than one band. It is derived with the SC lattice in consideration. The second method, the Möbius transformation method is far more general and very powerful.

3.2.1 Kalkstein & Soven method

Initially, a perfect infinite SC crystal is considered and separated artificially into two semi-infinite crystals along some imaginary cleavage plane in an arbitrary crystallographic direction [41, 42]. The perfect crystal has Hamiltonian H^0 and Green's function $G^0(\varepsilon - H^0 + i\delta) = 1$ where δ is a positive infinitesimal contribution to the energy, ε . The cleaved crystal has Hamiltonian H and Green's function $G(\varepsilon - H + i\delta) = 1$ where $H = H^0 + V$ and V is the potential arising due to the perturbation of the introduced surface. Thus $G(\varepsilon - H^0 - V + i\delta) = 1$ which implies $G(\varepsilon - H^0 - V + i\delta)G^0 = G^0$. Since G^0 is commutative the reformulation is possible $(\varepsilon - H^0 + i\delta)G^0 = 1$. Making this substitution above yields $G(1 - VG^0) = G^0$ and so the familiar Dyson's equation follows;

$$G = G^0 + GV G^0 \quad (22)$$

The coordinates of any atomic site in three dimensions are given by $\mathbf{R} = n_1 \mathbf{a}_1 + n_2 \mathbf{a}_2 + n_3 \mathbf{a}_3$ where $\mathbf{a}_1, \mathbf{a}_2, \mathbf{a}_3$ are vectors in the x, y, z -directions respectively. In this formalism however, the crystal is no longer infinite perpendicular to plane, thus the alternative basis $\mathbf{R} = n\mathbf{a} + \mathbf{R}_{\parallel}$ is used where $n = n_1$ and $\mathbf{a} = \mathbf{a}_1$ is a vector perpendicular to the cleaved surface and \mathbf{R}_{\parallel} refers to those sites parallel to the surface (surface

plane). The representation in the surface plane is Bloch-like so the transformation to the momentum space basis is as follows

$$|n, \mathbf{k}_{\parallel}\rangle = \frac{1}{\sqrt{N_{\parallel}}} \sum_{\mathbf{R}_{\parallel}} |n\mathbf{a} + \mathbf{R}_{\parallel}\rangle e^{i\mathbf{k}_{\parallel} \cdot \mathbf{R}_{\parallel}} \quad (23)$$

where N_{\parallel} is the total number of atomic sites in a plane parallel to the surface and \mathbf{k}_{\parallel} is the momentum space vector in this plane. The representation perpendicular to the surface plane is finite and thus the mixed Bloch-Wannier form is used

$$|\mathbf{k}\rangle = |\mathbf{k}_{\perp}, \mathbf{k}_{\parallel}\rangle = \frac{1}{\sqrt{N_{\perp}}} \sum_n |n\mathbf{a}, \mathbf{k}_{\parallel}\rangle e^{ink_{\perp}a} \quad (24)$$

where $N_{\perp}, \mathbf{k}_{\perp}$ are the total number of atomic sites, and momentum space vector respectively, perpendicular to plane. By the tight binding model, the Green's function of the unperturbed crystal is

$$\langle \mathbf{k} | G^0 | \hat{\mathbf{k}} \rangle = \frac{\delta_{\mathbf{k}-\hat{\mathbf{k}}}}{\varepsilon - E^0(\mathbf{k}) + i\delta} \quad (25)$$

where $\delta_{\mathbf{k}-\hat{\mathbf{k}}}$ is the Dirac-Delta function which is non-zero only for the very small range of values around $\mathbf{k} = \hat{\mathbf{k}}$. Here the result $\langle \mathbf{k} | H^0 | \hat{\mathbf{k}} \rangle = E^0(\mathbf{k})$ has been used, where

$$E^0(\mathbf{k}) = u + 2t[\cos(k_{\perp}a) + f(\mathbf{k}_{\parallel})] \quad (26)$$

where

$$f(\mathbf{k}_{\parallel}) = \cos(k_y a) + \cos(k_z a) \quad (27)$$

in the simple cubic case, as described here. Therefore

$$G^0(n - \acute{n}) = \left[\frac{1}{\sqrt{N_{\perp}}} \sum_{\acute{n}} \langle \mathbf{k}_{\parallel} | e^{-i\acute{n}k_{\perp}a} \right] \cdot \frac{1}{\varepsilon - E^0(\mathbf{k}) + i\delta} \cdot \left[\frac{1}{\sqrt{N_{\perp}}} \sum_n |\mathbf{k}_{\parallel}\rangle e^{ink_{\perp}a} \right] \quad (28)$$

which yields the Green's function for each layer $|n| = (n - \acute{n})$,

$$G^0(n) = \frac{1}{N_{\perp}} \sum_{k_{\perp}} \frac{e^{i|n|k_{\perp}a}}{\varepsilon - E^0(\mathbf{k}) + i\delta} \quad (29)$$

This is conveniently solved by converting the summation to an integral and solving by contour integration.

$$G^0(n) = \frac{a}{2\pi} \int_{-\frac{\pi}{a}}^{\frac{\pi}{a}} dk_{\perp} \frac{e^{i|n|k_{\perp}a}}{\varepsilon - E^0(\mathbf{k}) + i\delta} \quad (30)$$

where the range of integration is over the first Brillouin zone. Putting $\omega = \varepsilon - u - 2tf(\mathbf{k}_{\parallel})$ yields

$$G^0(n) = \frac{a}{2\pi} \int_{-\frac{\pi}{a}}^{\frac{\pi}{a}} dk_{\perp} \frac{e^{i|n|k_{\perp}a}}{\omega + i\delta - 2t \cos(k_{\perp}a)}. \quad (31)$$

The substitution $z = e^{ik_{\perp}a}$ is made, so that $\cos(k_{\perp}a) = \frac{1}{2}(e^{ik_{\perp}a} + e^{-ik_{\perp}a}) = \frac{1}{2}(z + \frac{1}{z})$. Also $\frac{dz}{dk_{\perp}} = iaz$, $dk_{\perp} = \frac{dz}{iaz}$, thus integrating over the unit circle C yields

$$G^0(n) = \frac{a}{2\pi} \oint_C \frac{dz}{iaz} \cdot \frac{z^{|n|}}{\omega + i\delta - tz - tz^{-1}} \quad (32)$$

now letting $\delta \rightarrow 0$, and rearranging

$$G^0(n) = \frac{i}{2\pi} \oint_C dz \frac{z^{|n|}}{t(z^2 - \omega z/t + 1)} \quad (33)$$

the denominator of the integrand is equivalent to

$$\left(z - \frac{\omega - \sqrt{\omega^2 - 4t^2}}{2t}\right) \left(z - \frac{\omega + \sqrt{\omega^2 - 4t^2}}{2t}\right) \quad (34)$$

thus the singularities are at $z = \frac{\omega - \sqrt{\omega^2 - 4t^2}}{2t}$ and $z = \frac{\omega + \sqrt{\omega^2 - 4t^2}}{2t}$ but only one of these is within the unit circle, and this depends on the sign and ratio of certain quantities. Thus there is one residue, and when $\omega > 0$, ($\omega^2 - 4t^2 \geq 0$) it is at the former, thus

$$G^0(n) = \frac{i}{2\pi} \left[\frac{2\pi i}{t} \cdot \frac{z^{|n|}}{z - \frac{\omega + \sqrt{\omega^2 - 4t^2}}{2t}} \right] \Big|_{z = \frac{\omega - \sqrt{\omega^2 - 4t^2}}{2t}} \quad (35)$$

$$G^0(n) = \frac{1}{\alpha} \left(\frac{\omega - \alpha}{2t} \right)^{|n|} \quad (36)$$

where $\alpha = \sqrt{\omega^2 - 4t^2}$. When $\omega < 0$, ($\omega^2 - 4t^2 \geq 0$) the residue is at $z = (\omega + \sqrt{\omega^2 - 4t^2})/2t$, thus

$$G^0(n) = \frac{1}{-\alpha} \left(\frac{\omega + \alpha}{2t} \right)^{|n|}. \quad (37)$$

These two expressions can be put more succinctly as follows

$$G^0(n) = \frac{1}{\alpha \text{sign}(\omega)} \left(\frac{\omega - \alpha \text{sign}(\omega)}{2t} \right)^{|n|} \quad (38)$$

where $|n|$ is used since symmetry is required on either side of the cleavage plane. The final situation to consider is $\omega^2 - 4t^2 < 0$, in which case the residue is at $z = (\omega + \sqrt{\omega^2 - 4t^2})/2t$, or $z = (\omega + i\sqrt{4t^2 - \omega^2})/2t$. Thus

$$G^0(n) = \frac{i}{\mu} \left(\frac{\omega + i\mu}{2t} \right)^{|n|} \quad (39)$$

where

$$\mu = \sqrt{4t^2 - \omega^2}. \quad (40)$$

Returning to Dyson's equation (22), which in this representation reduces to the set of equations

$$G(m, n) = G^0(m, n) + \sum_{p, q} G^0(m, p) V(p, q) G(q, n) \quad (41)$$

but $V(m, n)$ is only non-zero for sites on opposing surface planes so p and q only yield non-zero values when they are equal to -1 or 0 and not equal to each other. Thus in full

$$G(m, n) = G^0(m, n) + G^0(m, -1) V(-1, 0) G(0, n) + G^0(m, 0) V(0, -1) G(-1, n) \quad (42)$$

but $G(-1, n) = 0$ for $n \geq 0$ as these refer to sites on different crystals. Thus

$$G(m, n) = G^0(m, n) + G^0(m, -1)V(-1, 0)G(0, n) \quad (43)$$

and $V(-1, 0)$ must be equal to $-t$, since V exists to cancel the interactions between the exposed surfaces. Also since G^0 depends only upon the difference $m - n$ since it is the Green's function for the perfect crystal ^[42],

$$G(m, n) = G^0(m - n) - tG^0(m + 1)G(0, n) \quad (44)$$

Thus the semi-infinite Green's function for the cleaved crystal is

$$G(m, m) = G^0(0) - tG^0(m + 1)G(0, m) \quad (45)$$

where

$$G(0, m) = G^0(-m) - tG^0(1)G(0, m) \quad (46)$$

thus

$$[1 + tG^0(1)] G(0, m) = G^0(-m) \quad (47)$$

thus

$$G(m, m) = G^0(0) + [-tG^0(m + 1)] G^0(-m) [1 + tG^0(1)]^{-1}, \quad m \geq 0. \quad (48)$$

At this stage the Green's function may be elicited from equations (38) & (39). Starting with Equation (38), i.e. the case $\omega^2 - 4t^2 > 0$, where here the substitution $\eta = \alpha \text{sign}(\omega)$ is made

$$G(m, m) = \frac{1}{\eta} - \frac{t}{\eta^2} \left(\frac{\omega - \eta}{2t} \right)^{|2m+1|} \left[1 + \frac{1}{\eta} \left(\frac{\omega - \eta}{2} \right) \right]^{-1} \quad (49)$$

thus

$$G(m, m) = \frac{1}{\eta} + \left[\frac{1}{\eta} \left(\frac{\eta - \omega}{2} \right) \left(\frac{\omega - \eta}{2t} \right)^{2|m|} \right] \left[\frac{\eta + \omega}{2} \right]^{-1} \quad (50)$$

thus

$$G(m, m) = \frac{1}{\eta} \left[1 + \left(\frac{\omega - \eta}{2t} \right)^{2|m|} \frac{(\eta - \omega)}{(\eta + \omega)} \right], \quad \omega^2 - 4t^2 > 0 \quad (51)$$

And now the case $\omega^2 - 4t^2 < 0$, obtained by substituting Equation (39) into Equation (48)

$$G(m, m) = \frac{i}{\mu} + \left[-\frac{it}{\mu} \left(\frac{\omega + i\mu}{2t} \right)^{|m+1|} \right] \left[\frac{i}{\mu} \left(\frac{\omega + i\mu}{2t} \right)^{|m|} \right] \left[1 + \frac{it}{\mu} \left(\frac{\omega + i\mu}{2t} \right) \right]^{-1} \quad (52)$$

thus

$$G(m, m) = \frac{i}{\mu} + \left[-\frac{i}{\mu} \left(\frac{\omega + i\mu}{2} \right) \left(\frac{\omega + i\mu}{2t} \right)^{2|m|} \right] \left[\frac{-i\mu + \omega}{2} \right]^{-1} \quad (53)$$

thus

$$G(m, m) = \frac{i}{\mu} \left[1 + \left(\frac{\omega + i\mu}{2t} \right)^{2|m|} \frac{(\mu - i\omega)}{(\mu + i\omega)} \right], \quad \omega^2 - 4t^2 < 0 \quad (54)$$

Throughout the remainder of this communication the convention $G(0,0) \equiv g_0$ is used. To demonstrate one of the potential uses of SGFs, the local density of states (LDOS) will be calculated at different atomic layers. In the mixed representation, the LDOS is defined by

$$\rho_m(\varepsilon) = \frac{-1}{\pi N_{\parallel}} \sum_{\mathbf{k}_{\parallel}} \Im \{G(m, m, \mathbf{k}_{\parallel})\} \quad [42]. \quad (55)$$

Equation (55) is calculated numerically by converting the summation into a double integral over the two dimensional Brillouin zone

$$\rho_m(\varepsilon) = \frac{-a^2}{4\pi^3} \int_{-\frac{\pi}{a}}^{\frac{\pi}{a}} \int_{-\frac{\pi}{a}}^{\frac{\pi}{a}} dk_y dk_z G(m, m, k_y, k_z). \quad (56)$$

Local Density of States (LDOS)

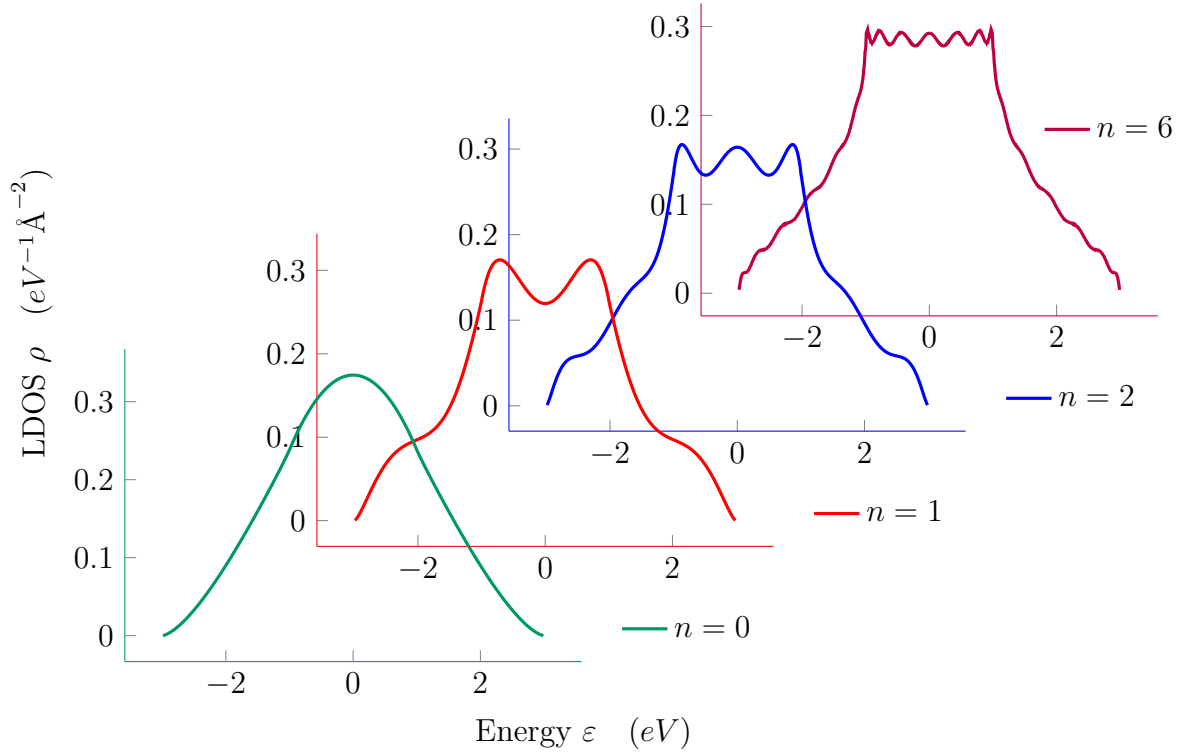


Figure 4: LDOS with potential $u = 0\text{eV}$, $t = 0.5\text{eV}$. The plots above are for atomic layers of various depth from the surface, n , with the surface layer $n = 0$. Results calculated using C++98 with Clang compiler.

Figure 4 shows the calculated LDOS against energy for up to the sixth atomic plane. The following results are noted; the profile of the LDOS for larger n heals quickly to the bulk value, and has two distinct zones.

The region $-2t \leq \varepsilon \leq 2t$, the width of the allowed values of the in-plane element of the dispersion relation Equation (27), is some periodic function that flattens out the larger the value n chosen. This region is flat in the case of an infinite crystal. The

region $\varepsilon < -2t, \varepsilon > 2t$ appears to be almost parabolic. The on-site energy u displaces the whole function along the ε axis. The full width of the function is $12t$, the width of the dispersion relation. By definition, integrating the LDOS gives 1.

3.2.2 Method of Adlayer

A simple method to gain SGFs for additional 'deposited' atomic planes is the so called 'Adlayer' method, which essentially connects the previous SGF to the deposited SGF through Dyson's equation. It is an iterative equation depending only on the previous SGF through each computation. To calculate this it can be noted that Equation (41) is equally valid for matrices, and so Equation (43) is rewritten in this context to find \mathbf{G}_{11} , the first Green's function deposited on the substrate

$$\mathbf{G}_{11}^1 = \mathbf{G}_{11}^0 + \mathbf{G}_{11}^1 \mathbf{V}_{11} \mathbf{G}_{11}^0 + \mathbf{G}_{10}^1 \mathbf{V}_{01} \mathbf{G}_{11}^0 \quad (57)$$

Using Equation (41) to find \mathbf{G}_{10}^1

$$\mathbf{G}_{10}^1 = \mathbf{G}_{11}^1 \mathbf{V}_{10} \mathbf{G}_{00}^0 \quad (58)$$

since $\mathbf{V}_{00} = \mathbf{0}$, $\mathbf{G}_{10}^0 = \mathbf{0}$. Substituting Equation (58) into Equation (57) yields

$$\mathbf{G}_{11}^1 = \mathbf{G}_{11}^0 + \mathbf{G}_{11}^1 \mathbf{V}_{11} \mathbf{G}_{11}^0 + \mathbf{G}_{11}^1 \mathbf{V}_{10} \mathbf{G}_{00}^0 \mathbf{V}_{01} \mathbf{G}_{11}^0 \quad (59)$$

which can be manipulated to

$$(\mathbf{G}_{11}^1)^{-1} \mathbf{G}_{11}^1 = (\mathbf{G}_{11}^1)^{-1} \mathbf{G}_{11}^0 + (\mathbf{G}_{11}^1)^{-1} \mathbf{G}_{11}^1 \mathbf{V}_{11} \mathbf{G}_{11}^0 + (\mathbf{G}_{11}^1)^{-1} \mathbf{G}_{11}^1 \mathbf{V}_{10} \mathbf{G}_{00}^0 \mathbf{V}_{01} \mathbf{G}_{11}^0 \quad (60)$$

which simplifies to

$$\mathbf{I} = (\mathbf{G}_{11}^1)^{-1} \mathbf{G}_{11}^0 + \mathbf{V}_{11} \mathbf{G}_{11}^0 + \mathbf{V}_{10} \mathbf{G}_{00}^0 \mathbf{V}_{01} \mathbf{G}_{11}^0 \quad (61)$$

which can be reordered to yield

$$\mathbf{I} = ((\mathbf{G}_{11}^1)^{-1} + \mathbf{V}_{11} + \mathbf{V}_{10} \mathbf{G}_{00}^0 \mathbf{V}_{01}) \mathbf{G}_{11}^0 \quad (62)$$

thus

$$(\mathbf{G}_{11}^0)^{-1} = (\mathbf{G}_{11}^1)^{-1} + \mathbf{V}_{11} + \mathbf{V}_{10} \mathbf{G}_{00}^0 \mathbf{V}_{01} \quad (63)$$

so upon reordering and taking appropriate inverses

$$\mathbf{G}_{11}^1 = ((\mathbf{G}_{11}^0)^{-1} - \mathbf{V}_{11} - \mathbf{V}_{10} \mathbf{G}_{00}^0 \mathbf{V}_{01})^{-1} \quad (64)$$

Now from the definition of Green's function, the term $\mathbf{G}_{11}^0 = (\boldsymbol{\varepsilon})^{-1}$. $\mathbf{V}_{11} = \mathbf{u}$, where \mathbf{u} is the potential of the adlayered plane, $\mathbf{V}_{10} = \mathbf{t}^\dagger$ is the hopping from the adlayer to the substrate and $\mathbf{V}_{01} = \mathbf{t}$ is the hopping from the substrate to the adlayer. Where the indices of the Green's functions are the same, this indicates they are SGFs and thus Equation (64) can be presented as follows

$$\mathbf{g}_1 = (\mathbf{v} - \mathbf{t}^\dagger \mathbf{g}_0 \mathbf{t})^{-1} \quad (65)$$

where $\mathbf{v} = \boldsymbol{\varepsilon} - \mathbf{u}$. Repeating the adlayer process leads to the iterative formula

$$\mathbf{g}_n = (\mathbf{v} - \mathbf{t}^\dagger \mathbf{g}_{n-1} \mathbf{t})^{-1} \quad (66)$$

an exact form of a closed-form SGF in terms of \mathbf{g}_0 . What makes Equation (66) particularly interesting is the fact that different materials may be stacked upon each other without loss of accuracy. Thus it makes a good vehicle for studying multilayer systems.

3.2.3 Möbius Transformation method

The Kalkstein & Soven method outlined in Subsection 3.2.1 is one of the first closed-form Green's function formalisms derived. However, it fails to describe the case when Green's functions are matrices. For any realistic calculations, this becomes a requirement. An alternative method known as the Möbius Transformation method, or the method of Umerski [43] is used in these cases.

As in Section 3.2.1, the formalism required is that of a crystal, infinite along two axes (in plane) and finite in the growth direction (out-of-plane). Due to the structure of the crystal here, the standard tight-binding model in three dimensions outlined in Section 3.1 does not apply. Instead, the Hamiltonian must be diagonalised in-plane, but unless the crystal is SC, in general more than one basis atom must be considered to cater for the principle layers, such that there is translational symmetry in the growth direction, and so that interactions between these layers are accounted for.

Thus the location of any atomic site is given by $\mathbf{R} = \mathbf{R}_{\parallel} + \mathbf{R}_{\perp} + \mathbf{s}_i$, where $\mathbf{R}_{\parallel} = n_1 \mathbf{a}_1 + n_2 \mathbf{a}_2$ represents the sites on the plane parallel to the surface per basis atom, with integers n_1, n_2 and $\mathbf{a}_1, \mathbf{a}_2$ are the primitive translation vectors. The vector $\mathbf{R}_{\perp} = n \mathbf{a}$ represents the sites along the growth direction, generally the y -axis (k_y -axis in reciprocal space) such that $\mathbf{a} = a(0, 1, 0)$. n is the n th plane from the SGF and a is the distance between planes. \mathbf{s}_i is the position vector of the basis atoms. For example, bcc has the two basis atoms at $\mathbf{s}_1 = (0, 0, 0), \mathbf{s}_2 = \frac{a}{2}(1, 1, 1)$, the two sites in a primitive (symmetrically irreducible) cell.

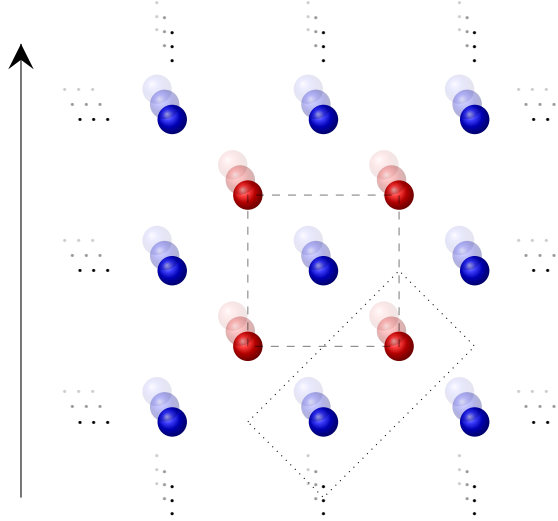


Figure 5: The above figure shows a bcc lattice with atomic planes stacked in the growth direction. The dashed region indicates the first Brillouin zone centred on the site at $(0,0,0)$, and the dotted region shows a primitive cell enclosing sites $a(0, -1, 0)$ and $\frac{a}{2}(1, -1, 1)$. It can be observed that the primitive cells tessellate completely, which forms part of their definition.

The representation in k -space is thus

$$|\mathbf{k}_{\parallel}, \mathbf{k}_{\perp}, \mathbf{s}_i\rangle = \frac{1}{\sqrt{N_{\parallel}}} \sum_{\mathbf{R}_{\parallel}=-\infty}^{\infty} e^{i\mathbf{k}_{\parallel} \cdot \mathbf{R}_{\parallel}} |\mathbf{R}_{\parallel}, \mathbf{R}_{\perp}, \mathbf{s}_i\rangle \quad (67)$$

where \mathbf{k}_{\parallel} is the in-plane momentum space vector, and \mathbf{k}_{\perp} is the momentum space vector in the growth direction (out of plane). N_{\parallel} is the number of atomic sites in-plane. Thus the diagonalised Hamiltonian can be given by

$$\langle \mathbf{k}_{\parallel}, \mathbf{k}_{\perp}, \mathbf{s}_i | \mathbf{H} | \mathbf{k}'_{\parallel}, \mathbf{k}'_{\perp}, \mathbf{s}_j \rangle = \frac{1}{N_{\parallel}} \sum_{\mathbf{R}_{\parallel} \mathbf{R}'_{\parallel}} e^{i(\mathbf{k}'_{\parallel} \cdot \mathbf{R}'_{\parallel} - \mathbf{k}_{\parallel} \cdot \mathbf{R}_{\parallel})} \langle \mathbf{R}_{\parallel}, \mathbf{R}_{\perp}, \mathbf{s}_i | \mathbf{H} | \mathbf{R}'_{\parallel}, \mathbf{R}'_{\perp}, \mathbf{s}_j \rangle \quad (68)$$

An example is now given, the bcc structure. In this case, the basis atoms are located at $\mathbf{s}_1 = (0, 0, 0)$, $\mathbf{s}_2 = \frac{a}{2}(1, 1, 1)$. The principle lattice vectors are $\mathbf{a}_1 = a(1, 0, 0)$, $\mathbf{a}_2 = a(0, 0, 1)$, $\mathbf{a}_3 = a(0, 1, 0)$, with $\mathbf{R}_{\parallel} = n_1 \mathbf{a}_1 + n_2 \mathbf{a}_2$, $\mathbf{R}_{\perp} = n_3 \mathbf{a}_3$. The potential of the on-site element of the first basis atom at \mathbf{s}_1 is \mathbf{u}_1 and the potential of the on-site element of the second basis atom at \mathbf{s}_2 is \mathbf{u}_2 , where the bold typeface of these potentials indicates they are themselves matrices in multi-band models. In this example, off-site terms up to second nearest neighbours are considered, where the potential induced by first nearest neighbour interactions is \mathbf{t}_1 , and second nearest neighbour interactions is \mathbf{t}_2 . Using Equation (68) in conjunction with Equation (6) yields

$$\langle \mathbf{k}_{\parallel}, n_3, \mathbf{s}_i | \mathbf{H} | \mathbf{k}'_{\parallel}, n_3, \mathbf{s}_j \rangle = \delta_{\mathbf{k}_{\parallel} \mathbf{k}'_{\parallel}} \sum_{\mathbf{d} \in N_{\parallel}} e^{i\mathbf{k}_{\parallel} \cdot \mathbf{d}} \langle \mathbf{0}, 0, \mathbf{s}_i | \mathbf{H} | \mathbf{d}, 0, \mathbf{s}_j \rangle \quad (69)$$

where \mathbf{d} is the set including the position vectors of the on-site potential, the nearest neighbours and the next nearest neighbours. The Hamiltonian diagonalised in-plane takes the following form

$$\mathbf{H} = \begin{pmatrix} \mathbf{u} & \mathbf{t} & \mathbf{0} & \mathbf{0} \\ \mathbf{t}^{\dagger} & \mathbf{u} & \ddots & \mathbf{0} \\ \mathbf{0} & \ddots & \ddots & \mathbf{t} \\ \mathbf{0} & \mathbf{0} & \mathbf{t}^{\dagger} & \mathbf{u} \end{pmatrix} \quad (70)$$

where \mathbf{u} , the on-plane potential, and \mathbf{t} , the out-of-plane hopping are $2N \times 2N$ matrices with each $N \times N$ element being dependent on the interaction of the two basis atoms, so that

$$\mathbf{u} = \begin{pmatrix} \mathbf{u}_{11} & \mathbf{u}_{12} \\ \mathbf{u}_{12}^{\dagger} & \mathbf{u}_{22} \end{pmatrix}, \quad \mathbf{t} = \begin{pmatrix} \mathbf{t}_{11} & \mathbf{t}_{12} \\ \mathbf{t}_{21} & \mathbf{t}_{22} \end{pmatrix} \quad (71)$$

where the diagonal elements of \mathbf{u} are given by

$$\begin{aligned} \mathbf{u}_{ii} &= \langle \mathbf{k}_{\parallel}, n_3, \mathbf{s}_i | \mathbf{H} | \mathbf{k}_{\parallel}, n_3, \mathbf{s}_i \rangle \\ &= \sum_{\mathbf{d} \in N_{\parallel}} e^{i\mathbf{k}_{\parallel} \cdot \mathbf{d}} \langle \mathbf{0}, 0, \mathbf{s}_i | \mathbf{H} | \mathbf{d}, 0, \mathbf{s}_i \rangle \\ &= \mathbf{u}_i + 2\mathbf{t}_2 (\cos(k_x a) + \cos(k_z a)) \end{aligned} \quad (72)$$

using the generalised result from 3.1.1. The off-diagonal elements of \mathbf{u} are given by

$$\begin{aligned}
\mathbf{u}_{12} &= \langle \mathbf{k}_{\parallel}, n_3, \mathbf{s}_1 | \mathbf{H} | \mathbf{k}_{\parallel}', n_3, \mathbf{s}_2 \rangle \\
&= \sum_{\mathbf{d} \in N_{\parallel}} e^{i\mathbf{k}_{\parallel} \cdot \mathbf{d}} \langle \mathbf{0}, 0, \mathbf{s}_1 | \mathbf{H} | \mathbf{d}, 0, \mathbf{s}_2 \rangle \\
&= \mathbf{t}_1 \left(e^{i\mathbf{k}_{\parallel} \cdot (0,0,0)} + e^{i\mathbf{k}_{\parallel} \cdot (-a,0,0)} + e^{i\mathbf{k}_{\parallel} \cdot (0,0,-a)} + e^{i\mathbf{k}_{\parallel} \cdot (-a,0,-a)} \right) \\
&= \mathbf{t}_1 \left(1 + e^{-ik_x a} + e^{-ik_z a} + e^{-ia(k_x + k_z)} \right)
\end{aligned} \tag{73}$$

The diagonal elements of the matrix \mathbf{t} represent hopping to a site of the same basis, the first and second bases respectively. The off-diagonal elements represent hopping between bases. As before, these terms are restricted to up to next nearest neighbour terms.

The elements $\mathbf{t}_{ii} = \mathbf{t}_2$. It is easy to see that the element $\mathbf{t}_{12} = \mathbf{0}$ since it cannot represent nearest neighbour, or next nearest neighbour interactions. The element \mathbf{t}_{21} has exactly the same potential as $\mathbf{u}_{12}^{\dagger}$, since the same interactions are considered. Thus all terms of the Hamiltonian diagonalised in-plane have been derived, which concludes the example.

The definition of an N -dimensional Möbius transformation is as follows, let

$$\mathbf{A} = \begin{pmatrix} \mathbf{a} & \mathbf{b} \\ \mathbf{c} & \mathbf{d} \end{pmatrix} \tag{74}$$

Where \mathbf{A} is a $2N \times 2N$ matrix, and $\mathbf{a}, \mathbf{b}, \mathbf{c}, \mathbf{d}, \mathbf{z}$ are $N \times N$ matrices, then

$$\mathbf{A} \bullet \mathbf{z} = (\mathbf{az} + \mathbf{b})(\mathbf{cz} + \mathbf{d})^{-1} \tag{75}$$

is the left hand Möbius transformation of \mathbf{z} by \mathbf{A} . This leads to the following result, derived from Dyson's equation

$$\mathbf{g}_{n+1} = \mathbf{X} \bullet \mathbf{g}_n \tag{76}$$

where

$$\mathbf{X} = \begin{pmatrix} \mathbf{0} & \mathbf{t}^{-1} \\ -\mathbf{t}^{\dagger} & \mathbf{vt}^{-1} \end{pmatrix} \tag{77}$$

Here, each element is an $N \times N$ matrix and $\mathbf{0}$ is a matrix of zeros. The formalism here assumes a pseudo one-dimensional crystal grown along an arbitrary axis with $\mathbf{v} = \boldsymbol{\varepsilon} - \mathbf{u}$, where $\boldsymbol{\varepsilon}$ is a matrix whose diagonal elements are the energy, ε , and off-diagonal elements are zero. Equation (77) in its present form and the results that follow are defined for left hand surfaces. For right-hand surfaces, values \mathbf{t} and \mathbf{t}^{\dagger} must be interchanged. To generalise to a three-dimensional crystal, an alternative basis must be used such as that defined at the beginning of this subsection.

It can be shown that Equation (75) is associative with respect to multiplication. A consequence of this is that

$$\mathbf{A} \bullet (\mathbf{B} \bullet \mathbf{z}) = (\mathbf{AB}) \bullet \mathbf{z} \tag{78}$$

utilising this result and iterating Equation (66), it can be shown that

$$\mathbf{g}_n = (\mathbf{X})^n \bullet \mathbf{g}_0 \tag{79}$$

Equation (79) can be expressed as

$$\mathbf{g}_n = (\mathbf{O}\mathbf{\Lambda}^n\mathbf{O}^{-1}) \bullet \mathbf{g}_0 \quad (80)$$

Here \mathbf{O} is the matrix of eigenvectors of \mathbf{X} , arranged into columns with associated eigenvalues λ ordered in terms of their magnitudes as follows

$$|\lambda_1| \leq |\lambda_2| \leq \dots \leq |\lambda_{2N}| \quad (81)$$

$\mathbf{\Lambda}$ is the diagonal eigenvalue matrix

$$\mathbf{\Lambda} \equiv \text{diag}(\lambda_1, \lambda_2, \dots, \lambda_{2N}) \quad (82)$$

such that

$$\mathbf{O}^{-1}\mathbf{X}\mathbf{O} = \mathbf{\Lambda} \quad (83)$$

Equation (80) has the remarkable property of demonstrating the existence of inter-atomic SGF representations due to the continuity at non-integer n .

To obtain the closed form SGF \mathbf{g}_0 for semi-infinite systems, the ordering of the eigenvalues must be more precise [25]. To do this, a small positive imaginary contribution is made to the energy, such that $\varepsilon \rightarrow \varepsilon + i\delta$ but to enforce $n \rightarrow \infty$ before $\delta \rightarrow 0$. This has the effect of shifting pairs of eigenvalues on the unit circle such that one moves just inside the circle and the other just outside. The former are in $\mathbf{\Lambda}_1$ and the latter are in $\mathbf{\Lambda}_2$, such that as $n \rightarrow \infty$, $\mathbf{\Lambda}_1^n, \mathbf{\Lambda}_2^{-n} \rightarrow 0$, where

$$\mathbf{\Lambda} = \begin{pmatrix} \mathbf{\Lambda}_1 & \mathbf{0} \\ \mathbf{0} & \mathbf{\Lambda}_2 \end{pmatrix} \quad (84)$$

The linearly transformed SGF is introduced

$$\mathbf{f}_0 = \mathbf{O}^{-1} \bullet \mathbf{g}_0 \quad (85)$$

whereupon combining with Equation (80)

$$\mathbf{g}_n = \mathbf{O} \bullet (\mathbf{\Lambda}_1^n \mathbf{f}_0 \mathbf{\Lambda}_2^{-n}) \quad (86)$$

as $n \rightarrow \infty$, $\mathbf{\Lambda}_1^n, \mathbf{\Lambda}_2^{-n} \rightarrow 0$. Thus the closed form SGF

$$\mathbf{g}_0 = \mathbf{O} \bullet \mathbf{0} \quad (87)$$

where $\mathbf{0}$ is the $N \times N$ null matrix.

4 Interlayer Exchange Coupling

In the absence of charge current in a magnetic multilayer, there remains a net flow of spincurrent between the two magnetic layers, provided the magnets are out of alignment. This is IEC, which oscillates between the FM and AF regimes as a function of spacer layer thickness [28]. Two means of calculating IEC are now given. The first is derived

from the spincurrent [25]. Using the Keldysh formalism [44], the spincurrent is essentially built from two components, $\langle j_{n-1} \rangle = \langle j_{n-1} \rangle_1 + \langle j_{n-1} \rangle_2$ [45]. with

$$\langle j_{n-1} \rangle_1 = \frac{1}{4\pi} \sum_{\mathbf{k}_{\parallel}} \int d\varepsilon \operatorname{Re} \operatorname{Tr} \{ (\mathbf{B} - \mathbf{A}) \boldsymbol{\sigma} \} [f(\varepsilon - \mu_L) + f(\varepsilon - \mu_R)] \quad (88)$$

the equilibrium term, and

$$\langle j_{n-1} \rangle_2 = \frac{1}{2\pi} \sum_{\mathbf{k}_{\parallel}} \int d\varepsilon \operatorname{Re} \operatorname{Tr} \left\{ \left[\mathbf{g}_L \mathbf{t} \mathbf{A} \mathbf{B} \mathbf{g}_R^{\dagger} \mathbf{t}^{\dagger} - \mathbf{A} \mathbf{B} + \frac{1}{2} (\mathbf{A} + \mathbf{B}) \right] \boldsymbol{\sigma} \right\} [f(\varepsilon - \mu_L) - f(\varepsilon - \mu_R)] \quad (89)$$

the transport term, named so since it is only non-zero when a bias is applied, where the electrochemical potentials in the right and left hand leads μ_L, μ_R are non-equal. Here $f(\varepsilon - \mu_{L(R)})$ is the Fermi-Dirac distribution function given by

$$f(\varepsilon - \mu_{L(R)}) = \frac{1}{1 + e^{(\varepsilon - \mu_{L(R)})/kT}} \quad (90)$$

The SGFs $\mathbf{g}_L, \mathbf{g}_R$, represent the left and right surfaces respectively, of a junction that is separated into two independent parts through some imaginary cleavage plane [46], see Fig 6. $\mathbf{A} = (\mathbf{I} - \mathbf{g}_R \mathbf{t}^{\dagger} \mathbf{g}_L \mathbf{t})^{-1}$, $\mathbf{B} = (\mathbf{I} - \mathbf{g}_R^{\dagger} \mathbf{t}^{\dagger} \mathbf{g}_L \mathbf{t})^{-1}$ and $\boldsymbol{\sigma}$ is one of the 2×2 Pauli matrices $\sigma_x, \sigma_y, \sigma_z$ whose choice is dependent on the spincurrent axis direction to be computed. As such, the out of plane spincurrent is required here, thus σ_y is selected for the computation. Equation (89) can instead yield the total charge current if $\frac{1}{2} \boldsymbol{\sigma}$ is replaced by $(e/\hbar) \mathbf{I}$ since $\langle j_{n-1} \rangle_1$ makes no contribution to charge current [47].

It is the equilibrium term (Equation (88)) that is of interest here, since it exists in the absence of applied bias and is responsible for the IEC.

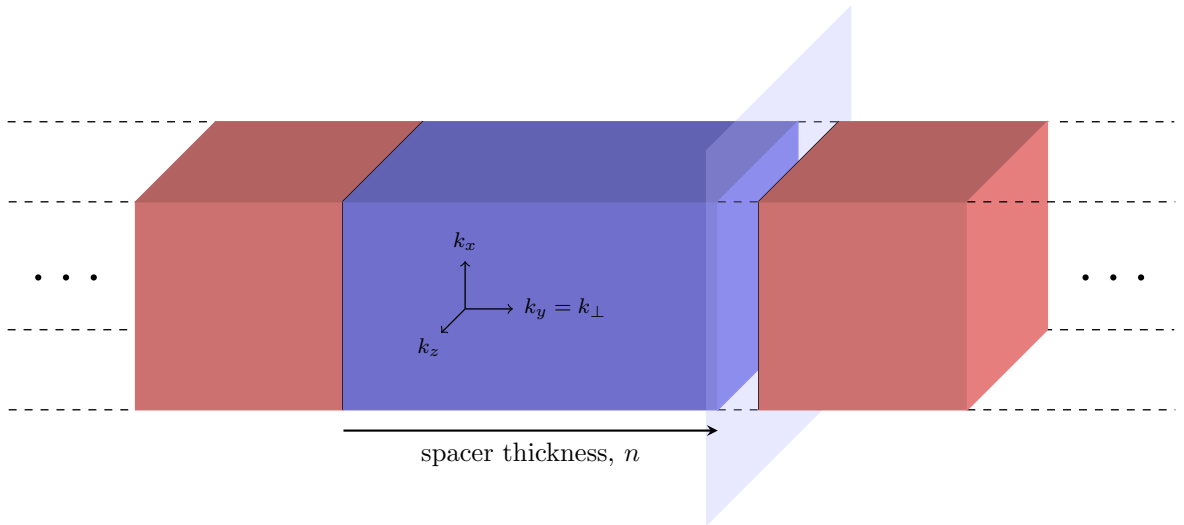


Figure 6: Schematic of a CPP trilayer. The red shaded areas are the semi-infinite leads, and the blue shaded area is the spacer. The spacer is grown from the LHS lead to thickness n , and for these calculations there is an imaginary cleavage plane passing between the n th spacer layer and the left surface of the RHS lead.

The system comprises (see Figure 6) of a polarising magnet (PM), a metallic spacer and a SM in its most minimal form. Often additional layers are required and in general, the system is grown with CPP if a charge current is passed through. Thus it may be modelled that the magnetisation of the PM lies in the xz plane at an angle θ to the z axis. The magnetisation of the SM is held parallel to the z axis. The closed-form SGFs of the leads may be calculated using Equation (87), the SGF of the spacer and SM are grown epitaxially using either Equation (66) or Equation (79). The author models \mathbf{g}_R as the SGF of the SM in contact with the spacer, and \mathbf{g}_L as the SGF of the outermost spacer layer, but without loss of generality, any such neighbouring SGFs may be used. The matrix of on-site potentials of the PM have been rotated by θ as follows [48]

$$\mathbf{u}_{1\theta} = \mathbf{s}^{-1} \mathbf{u}_{10} \mathbf{s} \quad (91)$$

where

$$\mathbf{s} = \begin{pmatrix} \cos(\frac{\theta}{2}) & \sin(\frac{\theta}{2}) \\ -\sin(\frac{\theta}{2}) & \cos(\frac{\theta}{2}) \end{pmatrix} \quad (92)$$

It will be of interest to calculate the IEC of a system in the case of an applied bias. Unfortunately the integrand is not smooth with respect to energy, ε , so at least at this stage the zero-bias (equilibrium) case is considered. At equilibrium, Equation (88) yields the IEC $J(n)$, by integrating with respect to θ , from 0 to π

$$J(n) = \int_{\theta} d\theta \langle j_{n-1} \rangle_1 = \frac{1}{2\pi} \sum_{\mathbf{k}_{\parallel}} \int d\varepsilon \int_0^{\pi} d\theta \operatorname{Re} \operatorname{Tr} \{ (\mathbf{B} - \mathbf{A}) \boldsymbol{\sigma} \} f(\varepsilon - \mu) \quad (93)$$

A second method for describing IEC is derived directly from the thermodynamic potentials, which in equilibrium form is as follows [25]

$$J(n) = \frac{1}{n_{\parallel}} \sum_{\mathbf{k}_{\parallel}} \operatorname{Im} \int_{-\infty}^{+\infty} d\varepsilon f(\varepsilon - \mu) F(\varepsilon, \mathbf{k}_{\parallel}, n) \quad (94)$$

Where

$$F = \frac{1}{\pi} \ln[\det(\mathbf{A}_{\text{FM}}^{-1} \mathbf{A}_{\text{AF}})] \quad (95)$$

Where \mathbf{A}_{FM} is in the FM configuration ($\theta = 0$), and \mathbf{A}_{AF} is in the AF configuration ($\theta = \pi$) [25]. For computational efficiency purposes, these $2N \times 2N$ matrices may be reduced to $N \times N$ matrices by separating \mathbf{A} into its spin up and spin down constituents $\mathbf{A} = \boldsymbol{\alpha}^{\uparrow} \boldsymbol{\alpha}^{\downarrow}$, thus $\boldsymbol{\alpha}^{\sigma} = (\mathbf{I} - \mathbf{g}_R^{\sigma} \mathbf{t}^{\dagger} \mathbf{g}_L^{\sigma} \mathbf{t})^{-1}$. In this communication, efforts are directed at solving Equation (94), which is again very numerically intensive if attempted to solve directly. Instead, two methods to solve Equation (94) are outlined here. The method of summing Matsubara frequencies is described initially [49], followed by an outline of the stationary phase approximation (SPA) [50].

4.1 Method of Matsubara

The Matsubara frequency method takes advantage of the fact that $\int_{-\infty}^{\infty} d\varepsilon$ can be integrated over the contour comprised of a semicircle whose diameter lies on the real line. This calculation is possible, since $F(\varepsilon, \mathbf{k}_{\parallel}, n)$ is analytic in this domain and $\mapsto 0$, faster than $|\varepsilon| \mapsto \infty$ [51], and the Fermi function has an infinite number of isolated

singularities at $\varepsilon = \mu \pm (2m + 1)kT\pi i$ for $m = 0, 1, 2, \dots$, none of which lie on the real line. kT is the Boltzmann constant multiplied by temperature. The semicircle lies in the upper half plane, and so all singularities at γi for $\gamma < 0$ are discarded. The residues of $f(\varepsilon - \mu)$ lie at

$$\lim_{\omega \rightarrow (2m+1)\pi i} \frac{1}{1 + e^\omega} = \quad (96)$$

$$\lim_{\omega \rightarrow (2m+1)\pi i} \frac{1}{\frac{d}{d\varepsilon} e^\omega} = -kT \quad (97)$$

where $\omega = (\varepsilon - \mu)/kT$. Thus

$$\int_{-\infty}^{+\infty} d\varepsilon f(\varepsilon - \mu) F(\varepsilon, \mathbf{k}_{\parallel}, n) \equiv 2\pi i \left(-kT \sum_{m=0}^{\infty} F(\mu + (2m + 1)kT\pi i, \mathbf{k}_{\parallel}, n) \right) \quad (98)$$

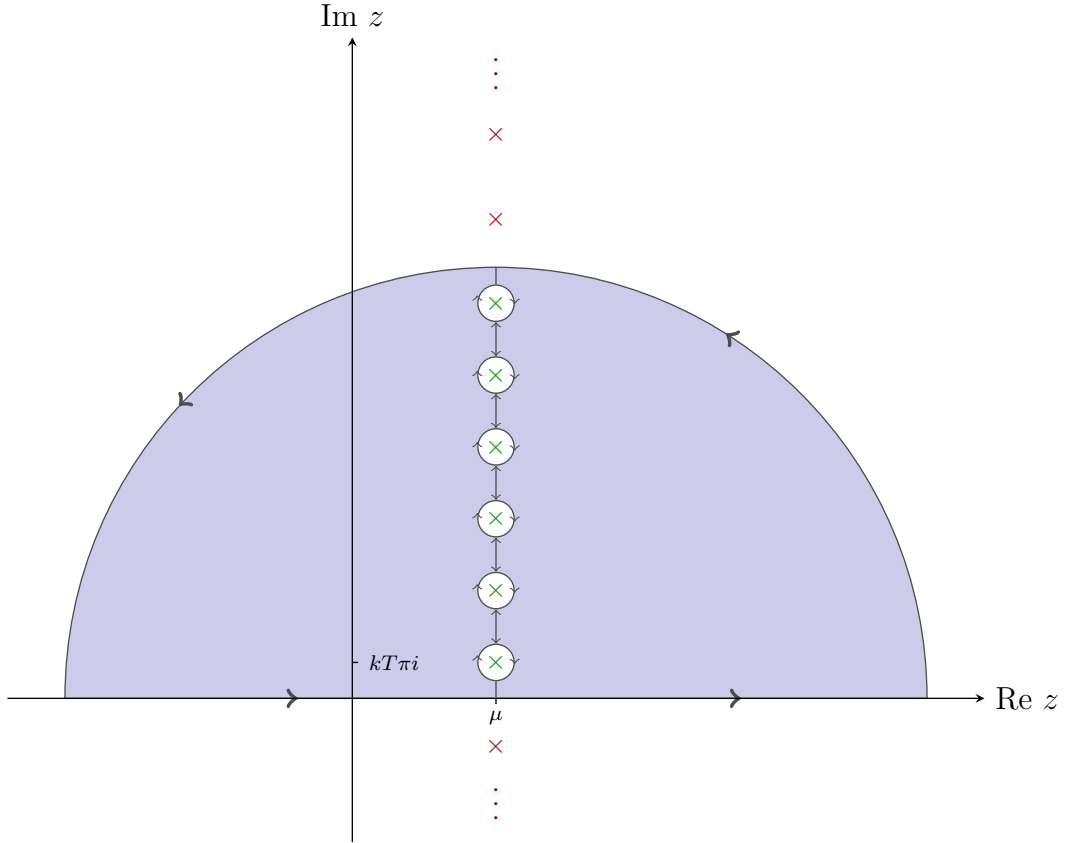


Figure 7: Summation of Matsubara frequencies. The integral $\int_{-\infty}^{+\infty}$ is replaced with a semi-circular contour of infinite radius. Since $f(\varepsilon - \mu)$ has residue $-kT$ at each of its poles $\mu + (2m + 1)kT\pi i$, for $m = 0, 1, \dots$, and F is analytic throughout the domain and converges for $\sum_{m=0}^{\infty} F(\mu + (2m + 1)kT\pi i, \mathbf{k}_{\parallel}, n)$, the integral may be replaced by the summation $2\pi i \left(-kT \sum_{m=0}^{\infty} F(\mu + (2m + 1)kT\pi i, \mathbf{k}_{\parallel}, n) \right)$. The summation converges rapidly getting very good agreement at nine terms. The contour pictured contains six of these poles.

The summation \mathbf{k}_{\parallel} in Equation (94) is over the two-dimensional Brillouin zone, which can be converted to an integral that can be evaluated numerically using quadra-

ture routines detailed in the Appendix [52], by making the following conversion

$$\frac{1}{n_{\parallel}} \sum_{\mathbf{k}_{\parallel}} \longrightarrow \frac{1}{A_{\text{BZ}}} \int_{\mathbf{k}_{\parallel}} d^2 \mathbf{k}_{\parallel} \quad (99)$$

where A_{BZ} is the area of the Brillouin zone. Thus, the final form Equation (94) takes in this setting is

$$J(n) = -\frac{2\pi kT}{A_{\text{BZ}}} \text{Re} \sum_{m=0}^{\infty} \int_{\mathbf{k}_{\parallel}} d^2 \mathbf{k}_{\parallel} \{ F(\mu + (2m+1)kT\pi i, \mathbf{k}_{\parallel}, n) \} \quad (100)$$

4.2 Stationary Phase Approximation

$F(\varepsilon, \mathbf{k}_{\parallel}, n)$ is a quasi-periodic function of the continuous thickness n [25], with periods given by $(2\pi/k_{\perp 1}, 2\pi/k_{\perp 2}, \dots, 2\pi/k_{\perp p})$ where $\{k_{\perp 1}, k_{\perp 2}, \dots, k_{\perp p}\}$ are the values of the perpendicular wave-vector k_{\perp} that satisfy the bulk dispersion relation, $E(\mathbf{k}_{\parallel}, k_{\perp}) = \varepsilon$. As such, $F(\varepsilon, \mathbf{k}_{\parallel}, n)$ may be expanded in a Fourier-type series (exactly a Fourier series in the one-band model)

$$F(\varepsilon, \mathbf{k}_{\parallel}, n) = \sum_{\mathbf{s}} c_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel}) e^{in\phi_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel})} \quad (101)$$

where $\phi_{\mathbf{s}} = (s_1 k_{\perp 1} + s_2 k_{\perp 2} + \dots + s_p k_{\perp p})$ with integers $\{s_1, s_2, \dots, s_p\}$ in \mathbf{s} , and $c_{\mathbf{s}}$ are the Fourier coefficients.

When this is substituted back into Equation (94), it can be observed that for large n , the imaginary component of the integrand oscillates rapidly as a function of \mathbf{k}_{\parallel} and ε . As such, the only non-zero contributions to the \mathbf{k}_{\parallel} integral come from the neighbourhood of points in k -space where $\phi_{\mathbf{s}}$ is stationary. Similarly, the only contributions to the ε integral that don't cancel are those that arise near the sharp cut-off that occurs at the chemical potential μ .

$c_{\mathbf{s}}$ is complex, so may be expressed in polar form $c_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel}) = |c_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel})| e^{i\psi_{\mathbf{s}}(\varepsilon)}$ thus

$$F(\varepsilon, \mathbf{k}_{\parallel}, n) = \sum_{\mathbf{s}} |c_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel})| e^{i(n\phi_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel}) + \psi_{\mathbf{s}}(\varepsilon))} \quad (102)$$

$\phi_{\mathbf{s}}$ is then expanded to second order in \mathbf{k}_{\parallel} about the stationary point, \mathbf{k}_{\parallel}^0

$$\phi_{\mathbf{s}}(\mathbf{k}_{\parallel}) = \phi_{\mathbf{s}}(\mathbf{k}_{\parallel}^0) + (\mathbf{k}_{\parallel} - \mathbf{k}_{\parallel}^0)^T \frac{\partial \phi_{\mathbf{s}}}{\partial \mathbf{k}_{\parallel}} + \frac{1}{2} (\mathbf{k}_{\parallel} - \mathbf{k}_{\parallel}^0)^T (\partial^2 \phi_{\mathbf{s}}) (\mathbf{k}_{\parallel} - \mathbf{k}_{\parallel}^0) + \dots \quad (103)$$

where $\partial \phi_{\mathbf{s}} / \partial \mathbf{k}_{\parallel}$, $(\partial^2 \phi_{\mathbf{s}})$ are 2×2 matrices, and $(\partial^2 \phi_{\mathbf{s}})$ is the Hessian matrix $(\partial^2 \phi_{\mathbf{s}})_{x,z} \equiv \partial^2 \phi_{\mathbf{s}} / \partial \mathbf{k}_{\parallel x} \partial \mathbf{k}_{\parallel z}$. The essence of the stationary phase approximation (SPA) is to assume the contributions are largest when the phase is stationary. Thus \mathbf{k}_{\parallel}^0 is the solution to $\partial \phi_{\mathbf{s}} / \partial \mathbf{k}_{\parallel} = 0$. Thus

$$J(n) = \frac{1}{A_{\text{BZ}}} \text{Im} \int_{-\infty}^{+\infty} d\varepsilon f(\varepsilon - \mu) \times \sum_{\mathbf{s}} |c_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel}^0)| e^{i(n\phi_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel}^0) + \psi_{\mathbf{s}}(\varepsilon))} \int_{\mathbf{k}_{\parallel}} d^2 \mathbf{k}_{\parallel} e^{i(\frac{n}{2} (\mathbf{k}_{\parallel} - \mathbf{k}_{\parallel}^0)^T (\partial^2 \phi_{\mathbf{s}}) (\mathbf{k}_{\parallel} - \mathbf{k}_{\parallel}^0))} \quad (104)$$

Where all terms not depending on \mathbf{k}_{\parallel} have been commuted to the left of the integral in k -space, and where the conversion, Equation (99) has been used. Using a standard matrix result from Gaussian integration [53], it can be shown that

$$\int_{-\infty}^{+\infty} \exp\left(\frac{i}{2}\mathbf{x}^T \mathbf{A} \mathbf{x}\right) d^2 \mathbf{x} \equiv \int_{-\infty}^{+\infty} \exp\left(\frac{i}{2} \sum_{j=1}^2 \lambda_j y_j^2\right) d^2 \mathbf{y} \quad (105)$$

where λ_j are the two eigenvalues of \mathbf{A} , and $\mathbf{y} = \mathbf{O}^T \mathbf{x}$ where \mathbf{O} is an orthogonal matrix obtained by diagonalising the symmetric matrix, \mathbf{A} such that $\mathbf{D} = \mathbf{O}^T \mathbf{A} \mathbf{O}$, where \mathbf{D} is a diagonal matrix.

Integrating with respect to each of the independent variables on the RHS of Equation (105) in turn, leads to the following result

$$\sqrt{\frac{2\pi i}{\lambda_1}} \sqrt{\frac{2\pi i}{\lambda_2}} \equiv \sqrt{\frac{2\pi}{|\lambda_1|}} e^{i\pi \operatorname{sgn}(\lambda_1)/4} \sqrt{\frac{2\pi}{|\lambda_2|}} e^{i\pi \operatorname{sgn}(\lambda_2)/4} \quad (106)$$

Equation (106) simplifies to

$$\begin{aligned} \frac{2\pi}{\sqrt{|\lambda_1 \lambda_2|}} e^{i\pi(\operatorname{sgn}(\lambda_1) + \operatorname{sgn}(\lambda_2))/4} &\equiv \frac{2\pi}{|\det(\mathbf{O}^T \mathbf{A} \mathbf{O})|^{\frac{1}{2}}} e^{i\pi(\operatorname{sgn}(\lambda_1) + \operatorname{sgn}(\lambda_2))/4} \\ &\equiv \frac{2\pi\tau}{|\det(\mathbf{O}^T \mathbf{A} \mathbf{O})|^{\frac{1}{2}}} \\ &\equiv \frac{2\pi\tau}{|\det(\mathbf{A})|^{\frac{1}{2}}} \end{aligned} \quad (107)$$

where $\tau = i$ if both eigenvalues of \mathbf{A} are positive, $\tau = -i$ if both eigenvalues are negative and $\tau = 1$ when both eigenvalues differ in sign. The final result in Equation (107) holds since $\det(\mathbf{O}) = 1$. The RHS integral in Equation (104) can be converted into the form of the RHS of Equation (105)

$$\int_{-\infty}^{+\infty} d^2 \mathbf{k}_{\parallel} e^{i(\frac{n}{2}(\mathbf{k}_{\parallel} - \mathbf{k}_{\parallel}^0)^T (\partial^2 \phi_{\mathbf{s}})(\mathbf{k}_{\parallel} - \mathbf{k}_{\parallel}^0))} \equiv \int_{-\infty}^{+\infty} d^2 \mathbf{y} e^{(\frac{i}{2} \mathbf{y}^T (n \partial^2 \phi_{\mathbf{s}}) \mathbf{y})} \quad (108)$$

since scalars are commutative, and upon the substitution $(\mathbf{k}_{\parallel} - \mathbf{k}_{\parallel}^0) = \mathbf{y}$ with $d^2 \mathbf{k}_{\parallel} = d^2(\mathbf{y} + \mathbf{k}_{\parallel}^0) = d^2 \mathbf{y}$. Thus Equation (108) yields

$$\frac{2\pi\tau}{n|\det(\partial^2 \phi_{\mathbf{s}})|^{\frac{1}{2}}} \quad (109)$$

The IEC presently takes the form

$$J(n) = \frac{2\pi}{nA_{\text{BZ}}} \operatorname{Im} \int_{-\infty}^{+\infty} d\varepsilon f(\varepsilon - \mu) \sum_{\mathbf{k}_{\parallel}^0} \sum_{\mathbf{s}} \frac{\tau |c_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel}^0)| e^{i(n\phi_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel}^0) + \psi_{\mathbf{s}}(\varepsilon))}}{|\det(\partial^2 \phi_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel}))|^{\frac{1}{2}}} \quad (110)$$

Next, $\psi_{\mathbf{s}}$ and $\phi_{\mathbf{s}}$ are expanded to first order about μ

$$\psi_{\mathbf{s}}(\varepsilon) \approx \psi_{\mathbf{s}}(\mu) + (\varepsilon - \mu)\psi'_{\mathbf{s}}(\mu), \quad \phi_{\mathbf{s}}(\varepsilon) \approx \phi_{\mathbf{s}}(\mu) + (\varepsilon - \mu)\phi'_{\mathbf{s}}(\mu) \quad (111)$$

where all dashed terms are differentiated with respect to ε . Making this substitution and bringing all terms that do not depend on ε to the left of the integral yields

$$J(n) = \frac{2\pi}{nA_{\text{BZ}}} \sum_{\mathbf{k}_{\parallel}^0} \sum_{\mathbf{s}} \text{Im} \left\{ \tau e^{i(n\phi_{\mathbf{s}}(\mu, \mathbf{k}_{\parallel}^0) + \psi_{\mathbf{s}}(\mu))} \times \int_{-\infty}^{+\infty} d\varepsilon f(\varepsilon - \mu) \frac{|c_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel}^0)| e^{i(n(\varepsilon - \mu)\phi'_{\mathbf{s}}(\mu, \mathbf{k}_{\parallel}^0) + (\varepsilon - \mu)\psi'_{\mathbf{s}}(\mu))}}{\left| \det \left(\partial^2 \phi_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel}^0) \right) \right|^{\frac{1}{2}}} \right\} \quad (112)$$

Next the substitution $x = (\varepsilon - \mu)/kT$ is made, thus $dx = d\varepsilon/kT$ and $f(\varepsilon - \mu) = \frac{1}{1+e^x}$ thus the integral takes the form

$$kT \int_{-\infty}^{+\infty} dx \frac{|c_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel}^0)|}{\left| \det \left(\partial^2 \phi_{\mathbf{s}}(\varepsilon, \mathbf{k}_{\parallel}^0) \right) \right|^{\frac{1}{2}}} \frac{e^{ixkT(n\phi'_{\mathbf{s}}(\mu, \mathbf{k}_{\parallel}^0) + \psi'_{\mathbf{s}}(\mu))}}{1 + e^x} \quad (113)$$

whereupon inspection of the left hand fraction within the integrand, given that it varies slowly with respect to ε and lies entirely on the real line, ε may be replaced by its first order term

$$kT \frac{|c_{\mathbf{s}}(\mu, \mathbf{k}_{\parallel}^0)|}{\left| \det \left(\partial^2 \phi_{\mathbf{s}}(\mu, \mathbf{k}_{\parallel}^0) \right) \right|^{\frac{1}{2}}} \int_{-\infty}^{+\infty} dx \frac{e^{ixkT(n\phi'_{\mathbf{s}}(\mu, \mathbf{k}_{\parallel}^0) + \psi'_{\mathbf{s}}(\mu))}}{1 + e^x} \quad (114)$$

using contour integration it can be shown that [\[51\]](#)

$$\int_{-\infty}^{+\infty} dx \frac{e^{ax}}{1 + e^x} = \frac{\pi}{\sin(\pi a)} \quad (115)$$

Thus Equation (114) becomes

$$\frac{kT |c_{\mathbf{s}}(\mu, \mathbf{k}_{\parallel}^0)| \pi}{\left| \det \left(\partial^2 \phi_{\mathbf{s}}(\mu, \mathbf{k}_{\parallel}^0) \right) \right|^{\frac{1}{2}} \sin \left(\pi i \left[n\phi'_{\mathbf{s}}(\mu, \mathbf{k}_{\parallel}^0) + \psi'_{\mathbf{s}}(\mu) \right] kT \right)} \equiv \frac{-ikT \pi |c_{\mathbf{s}}(\mu, \mathbf{k}_{\parallel}^0)|}{\left| \det \left(\partial^2 \phi_{\mathbf{s}}(\mu, \mathbf{k}_{\parallel}^0) \right) \right|^{\frac{1}{2}} \sinh \left(\pi kT \left[n\phi'_{\mathbf{s}}(\mu, \mathbf{k}_{\parallel}^0) + \psi'_{\mathbf{s}}(\mu) \right] \right)} \quad (116)$$

thus in full, the Stationary Phase Approximation of the IEC is

$$J(n) = -\frac{(2\pi)^2 kT}{2nA_{\text{BZ}}} \sum_{\mathbf{s}} \sum_{\mathbf{k}_{\parallel}^0} \text{Re} \left(\frac{\tau c_{\mathbf{s}} e^{in\phi_{\mathbf{s}}}}{\left| \det(\partial^2 \phi_{\mathbf{s}}) \right|^{\frac{1}{2}} \sinh(\pi kT [n\phi'_{\mathbf{s}} + \psi'_{\mathbf{s}}])} \right) \Big|_{\mu, \mathbf{k}_{\parallel}^0} \quad [54] \quad (117)$$

4.2.1 1-band Simple Cubic model

Given that in this simple case the Fermi surface has a single sheet, the period P is equal to π/k_{\perp} , and $\phi_s = 2sk_{\perp}$ [25]. k_{\perp} is obtained by rearranging the dispersion relation

$$\varepsilon = u + 2t (\cos(k_x a) + \cos(k_z a) + \cos(k_{\perp} a)) \quad (118)$$

where $\mathbf{k}_{\parallel} = (k_x, k_z)$. The summation over \mathbf{k}_{\parallel} is reduced to the sum of those values at which \mathbf{k}_{\parallel} is stationary. These values occur at $\partial\phi_s/\partial k_x = \partial\phi_s/\partial k_z = 0$. Since

$$k_{\perp} = \frac{1}{a} \cos^{-1} \left(\frac{\varepsilon - u}{2t} - \cos(k_x a) - \cos(k_z a) \right) \quad (119)$$

then

$$\frac{\partial k_{\perp}}{\partial k_x} = - \frac{\sin(k_x a)}{\sqrt{1 - \left(\frac{\varepsilon - u}{2t} - \cos(k_x a) - \cos(k_z a) \right)^2}} \quad (120)$$

and

$$\frac{\partial k_{\perp}}{\partial k_z} = - \frac{\sin(k_z a)}{\sqrt{1 - \left(\frac{\varepsilon - u}{2t} - \cos(k_x a) - \cos(k_z a) \right)^2}} \quad (121)$$

Thus there is one stationary point, at $\mathbf{k}_{\parallel} = \mathbf{0}$. The Fourier term

$$c_s = \frac{1}{P} \int_{n_0}^{n_0+P} dn F(n) e^{-i \frac{2\pi s n}{P}} \quad (122)$$

is integrated by utilising Equation (80), since this method is capable of handling non-integer n . The choice of n_0 may be arbitrary. The Hessian matrix, $(\partial^2 \phi_s)$ is evaluated at the stationary point \mathbf{k}_{\parallel}^0 . With these simplifications, the SPA may be restated as follows

$$J(n) = -\frac{kT}{2n} \sum_s \sum_{\mathbf{k}_{\parallel}^0} \text{Re} \left(\frac{\tau c_s e^{2ins k_{\perp}}}{2s |\det(\partial^2 k_{\perp})|^{1/2} \sinh(\pi kT [2ns k'_{\perp} + \psi'_s])} \right) \Big|_{\mu, \mathbf{k}_{\parallel}^0} \quad (123)$$

where

$$k'_{\perp} = \left(\frac{\partial \varepsilon}{\partial k_{\perp}} \right)^{-1} = (-2ta \sin(k_{\perp}))^{-1} \quad (124)$$

and ψ'_s is calculated numerically by comparing points on the gradient line over an infinitesimal distance such that the gradient is linear over these points. The Hessian matrix evaluates in this case to

$$(\partial^2 \phi_s)_{x,z} = \begin{pmatrix} -\frac{a \cos(k_x a)}{\sqrt{-\cos^2(k_{\perp} a) + 1}} + \frac{a \sin^2(k_x a) \cos(k_{\perp} a)}{(-\cos^2(k_{\perp} a) + 1)^{1/2}} & \frac{a \sin(k_x a) \sin(k_z a) \cos(k_{\perp} a)}{(-\cos^2(k_{\perp} a))^{1/2}} \\ \frac{a \sin(k_x a) \sin(k_z a) \cos(k_{\perp} a)}{(-\cos^2(k_{\perp} a))^{1/2}} & -\frac{a \cos(k_z a)}{\sqrt{-\cos^2(k_{\perp} a) + 1}} + \frac{a \sin^2(k_z a) \cos(k_{\perp} a)}{(-\cos^2(k_{\perp} a) + 1)^{1/2}} \end{pmatrix} \quad (125)$$

which is further simplified since $k_x, k_z = 0$. The results comparing the method of Matsubara, and SPA are presented in Figure 8.

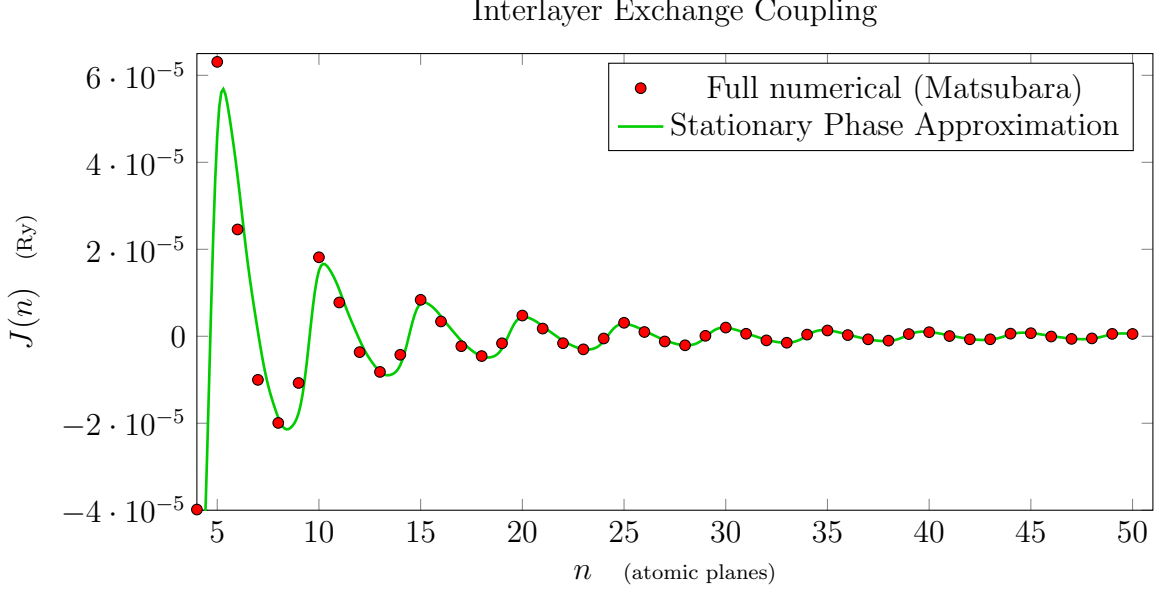


Figure 8: Interlayer Exchange coupling given by both the Matsubara sum method, where the summation of Matsubara frequencies is taken as \sum_0^9 , and SPA, where the summation of Fourier terms is taken over $\sum_{s=-5}^{s=5}$ though it can be observed that $c_0 = 0$. The potential in the leads $= -2.309Ry$ in the majority spin band and $= -0.95Ry$ in the minority band. The potential in the metallic spacer is $= -2.8Ry$. $t = 0.5Ry$, $\mu = (0 + 1 \times 10^{-5}i)Ry$ for SPA, and $(0 + 1 \times 10^{-10}i)Ry$ for Matsubara. $kT = 1.9 \times 10^{-3}Ry$ (room temperature). Results calculated using C++11 with Clang compiler, utilising the Eigen library. Concurrency provided through OpenMP.

4.3 Phase of Interlayer Exchange Coupling

It is known that IEC oscillates as a function of spacer thickness, but it is not so well understood what governs the phase of these oscillations, or the potential commercial applications that can take advantage of manipulating the phase. As discussed in Section 2.1, manipulating this phase may be used as a mechanism to induce switching of the SM making the system a good candidate for MRAM. In order to aid understanding of the origin of the phase, the SPA was analysed.

By perturbing quantities during modelling, and isolating terms during certain summations within Equation (117) in the one-band model, it became apparent the phase was dominated by the first Fourier term (using direct complex Fourier series expansion against SGFs \mathbf{g} , this term is c_{-1}), in particular since this term is complex, it is the phase ψ_s of the Fourier term that is of importance.

In order to understand the origin of this term more fully, it would be beneficial to negate the need to integrate over the spacer period, which is required using standard Fourier series expansion. There are two alternative methods to obtain the Fourier components. The first involves taking an inverse FFT. The second method replaces the SGF with the Möbius-transformed Green's function (MTGF)

$$\mathbf{f}_n = \mathbf{O}^{-1} \bullet \mathbf{g}_n \quad (126)$$

Now recall c_s are the Fourier components of F , Equation (95) with

$$\mathbf{A}^{\alpha\beta} = (\mathbf{I} - \mathbf{g}_R^\alpha \mathbf{t}^\dagger \mathbf{g}_L^\beta)^{-1} \quad (127)$$

where α, β are the spin directions, so that

$$F(\mathbf{n}) = F^{\uparrow\uparrow}(\mathbf{n}) + F^{\downarrow\downarrow}(\mathbf{n}) - F^{\uparrow\downarrow}(\mathbf{n}) - F^{\downarrow\uparrow}(\mathbf{n}) \quad (128)$$

with

$$F^{\alpha\beta} = \ln \det(\mathbf{A}^{\alpha\beta})^{-1} \quad (129)$$

and it can be shown that [\[55\]](#)

$$F^{\alpha\beta} = \text{Tr} [\ln (\mathbf{I} - \mathbf{f}_R^\alpha(0) \mathbf{f}_L^\beta(\mathbf{n}))] + \text{terms that cancel in Equation (128)} \quad (130)$$

where since $|\mathbf{f}| < 1$, Equation (130) may be expanded as a Taylor series

$$F^{\alpha\beta} \approx -\text{Tr} \left[\mathbf{f}_R^\alpha(0) \mathbf{f}_L^\beta(\mathbf{n}) + \frac{1}{2} (\mathbf{f}_R^\alpha(0))^2 (\mathbf{f}_L^\beta(\mathbf{n}))^2 + \frac{1}{3} (\mathbf{f}_R^\alpha(0))^3 (\mathbf{f}_L^\beta(\mathbf{n}))^3 + \dots \right] \quad (131)$$

but in the one band tight binding model, using Equations (86), (126), it is easy to see that $\mathbf{f}_L^\beta(\mathbf{n}) = \mathbf{f}_L^\beta(0) e^{-2ik_\perp \mathbf{n}}$, since $\lambda_1 = e^{-ik_\perp}$, $\lambda_2 = e^{ik_\perp}$. Thus Equation (131) can be re-expressed as follows

$$F^{\alpha\beta} \approx - \sum_s \frac{1}{s} (\mathbf{f}_R^\alpha \mathbf{f}_L^\beta) e^{-2isk_\perp \mathbf{n}} \quad (132)$$

where the independent variable is now dropped as the surface MTGFs are considered from here on.

Thus the MTGFs give a remarkably simple representation of the Fourier coefficients, c_s , giving a good platform to study the phase analytically.

It is the first term that dominates the phase of IEC, c_1

$$c_1^{\alpha\beta} = (f_R^\alpha f_L^\beta) \quad (133)$$

Continuing in the one band model and exploring an arbitrary MTGF, with Equation (126) yields

$$f = - \frac{(\lambda_2 t g - 1)}{(\lambda_1 t g - 1)} \quad (134)$$

and utilising Equation (87) together with knowledge that

$$\mathbf{O} = \begin{pmatrix} \mathbf{o}_1 & \mathbf{o}_2 \\ \boldsymbol{\tau} \mathbf{o}_1 \boldsymbol{\Lambda}_1 & \boldsymbol{\tau} \mathbf{o}_2 \boldsymbol{\Lambda}_2 \end{pmatrix} \quad (135)$$

with $\mathbf{o}_1, \mathbf{o}_2$ denoting the upper-left and upper-right sub-matrices of \mathbf{O} , and $\boldsymbol{\tau}$ is the hopping matrix of the lead, it is straightforward to show that $g = \tau^{-1} \lambda_1$ with λ_1 being the eigenvalue of the lead. This yields the following depending on whether the lead is an insulator or a conductor. In the case of a conductor

$$f = - \frac{(\frac{t}{\tau} e^{i(k_\perp - \kappa_\perp)} - 1)}{(\frac{t}{\tau} e^{-i(k_\perp + \kappa_\perp)} - 1)} \quad (136)$$

and in the case of an insulator

$$f = -\frac{\left(\frac{t}{\tau}e^{-\kappa_{\perp}}e^{ik_{\perp}} - 1\right)}{\left(\frac{t}{\tau}e^{-\kappa_{\perp}}e^{-ik_{\perp}} - 1\right)} \quad (137)$$

Where the Greek characters refer to quantities of the leads, and the Roman characters refer to quantities of the spacer, with $k_{\perp}, \kappa_{\perp}$ the perpendicular wave-vectors, and t, τ the one band hoppings.

Complex analysis shows the system of most interest here is that of the insulator since the numerator is the complex conjugate of the denominator of $-f$, thus $|f| = 1$, from which it follows that $|c_1| = 1$ establishing the dominance of this term. For the conducting spacer, $|f| < 1$ since the sign prefix before κ_{\perp} always ensures the magnitude of the numerator is less than that of the denominator.

Similarly, with κ_{\perp} being coupled to i in the conductor, the phase of the numerator of $-f$ will always be more advanced than that of the denominator. In the insulator, the two phases are equal and the total phase of the system is the sum of the phases of the denominator and numerator. It is thus expected the larger phase changes will be observed in an insulator, outlined in Fig. 9.

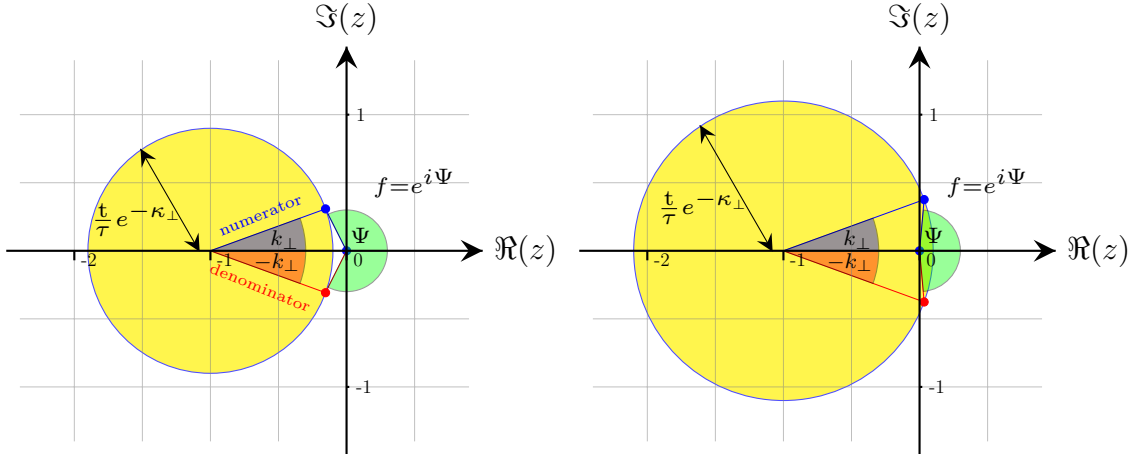


Figure 9: Depiction of Equation (137), the MTGF in the case of an insulating spacer. In each of the drawings, the blue point on the yellow disc indicates the numerator and the red point indicates the denominator. The blue shaded arc is the phase k_{\perp} , and the red shaded arc is the phase $-k_{\perp}$. The phase Ψ is of interest here, which is the phase of the MTGF f , depicted by the green shaded arc. The objective is to attempt to observe a phase change of Ψ from 0 to 2π . It is apparent that a large phase change of Ψ can be observed as the radius of the yellow disc, $\frac{t}{\tau}e^{-\kappa_{\perp}} = r$ moves from $r < 1$ to $r \geq 1$, when the phase k_{\perp} is small.

With this theory in account, a one dimensional SC one band model was used to compute the IEC at different ratios of hopping in the leads compared to the spacer. The results are shown in Fig. 10. The system explored has a weak insulating spacer so that the oscillations of the IEC were only slightly damped. Fig. 10 shows a strong phase dependence on the hopping ratios.

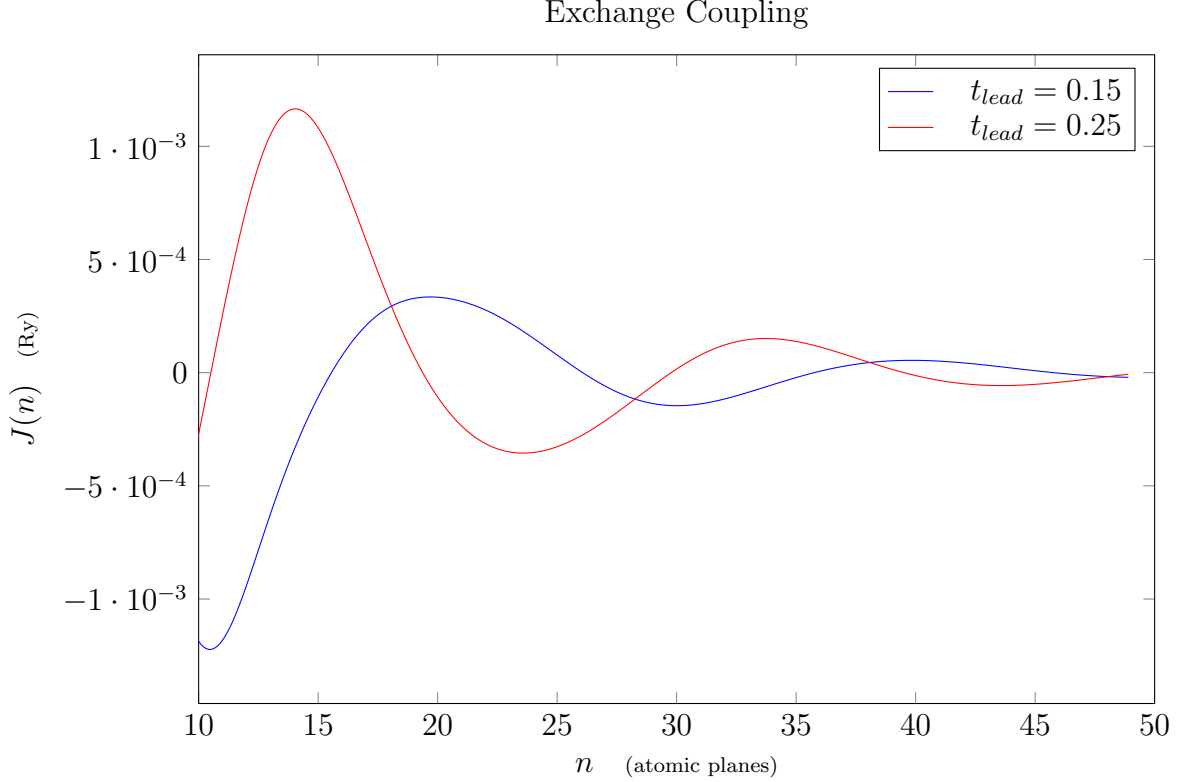


Figure 10: This figure shows the IEC in the one dimensional case, as a function of spacer thickness for two different hoppings in the leads in Rydbergs, showing the SPA calculations. The potential in the leads $= -0.7Ry$ in the majority spin band and $= -0.5Ry$ in the minority band. The potential in the spacer is $= -0.79Ry$, the hopping in the spacer is $t = 0.4Ry$, $\mu = (0 + 1 \times 10^{-5}i)Ry$. $kT = 1.9 \times 10^{-3}Ry$ (room temperature). Results calculated using C++11 with Clang compiler, utilising the Eigen library. Concurrency provided through OpenMP.

With the potentials used in Fig. 10, firstly considering the case when the hopping in the leads is $0.15Ry$, the leads behave as insulators in both spin bands and in each case the argument $r = \frac{t}{\tau}e^{-\kappa_{\perp}} < 1$ (see Fig. 9 for details).

Now considering the case when the hopping in the leads is $0.25Ry$, the leads behave as insulators in the majority spin band with $r < 1$. In the minority spin band, the leads have in the strictest sense become conducting, but this case is in fact the boundary term where Equations (136), (137) are equal. Here $r > 1$.

It is since in the minority spin band the radial term in f changes from $r \approx 0.89$ to $r = 1.60$, that the phase change in Fig. 10 is observed.

5 Conclusion & Future Developments

This communication has outlined the method of tight-binding and the Green's function formalism and applied these to describe the IEC. The IEC was given in its general form, then an analytical form was derived. Analysis of the analytic form (SPA) led to qualitative descriptions of the governance of phase in IEC.

The discovery by Umerski ^[37] regarding the 2π topologically protected phase shift has opened up an exciting area of research, with potential application to low energy MRAM devices. This is an end-goal that has been under active research for some time, though effort has been directed to STT-MRAM, which at this stage requires a large current density to enable switching. MRAM by virtue of its design, is non-volatile which is also of practical significance.

Future developments will include comprehending and making use of realistic tight-binding models with as many as nine bands, to make fully realistic calculations of systems to explore the phase dependence of the IEC. In order to simulate 'realness', surface roughness will be modelled and a thin switching layer will be considered. Different realistic structures will be modelled to explore the suitability of certain materials, though due to limitations of theory, these will of course be crystalline.

Due to the required gate voltage, the Matsubara frequency method will no longer be valid as a means of calculating IEC. Thus a full numerical computation requires direct integration over energy and two-dimensional k -space. This is very numerically intensive and the integrand as a function of energy is not smooth. To make this problem more tractable, the author has been developing a quadrature routine that is doubly adaptive and takes advantage of Brillouin zone symmetries. More details are given in the Appendix.

A Numerical Integration

A.1 Cunningham's Special Points set

In his 1974 paper ^[52], Cunningham describes the special points set that minimises the required function calls for converged integration of the two dimensional Brillouin zone (2DBZ), for a number geometries. His work follows that of Chadi and Cohen ^[56], where the integrable function is expanded as a Fourier series

$$f(\mathbf{k}) = f_0 + \sum_{m=1}^{\infty} f_m A_m(\mathbf{k}) \quad (138)$$

with

$$A_m(\mathbf{k}) = \sum_{|\mathbf{R}|=c_m} e^{i\mathbf{k}\cdot\mathbf{R}}, \quad m = 1, 2, \dots \quad (139)$$

The exact average over the 2DBZ is equal to f_0 , since the second term in Equation (138) vanishes. As such, upon imposing the restrictions

$$\sum_i \alpha_i A_m(\mathbf{k}) = 0, \quad m = 1, 2, \dots, N \quad (140)$$

$$\sum_i \alpha_i = 1 \quad (141)$$

where N is the number of functions $A_m(\mathbf{k})$ that satisfy Equation (140).

They found the most efficient means of integrating over the 2DBZ was by

$$f_0 = \sum_i \alpha_i f(\mathbf{k}_i) - \sum_{m>N}^{\infty} f_m \sum_i \alpha_i A_m(\mathbf{k}_i) \quad (142)$$

Though in primitive-rectangular or square lattices, a set of generating wave vectors could always be found that zeroes all functions $A_m(\mathbf{k})$, so the algorithm reduces to a 'mesh' of equally spaced points, displaced from the origin by $a(\frac{1}{2}, \frac{1}{2})$ where a is the distance between the neighbouring points.

For these simple lattice types, the integral is performed as follows;

$$\iint_{\text{BZ}} f(k_x, k_z) dk_x dk_z = F_N \quad (143)$$

where in the case of the square lattice

$$F_N = \frac{8A_{\text{BZ}}}{N^2} \sum_{n=0}^{2N} \left(\sum_{m=0}^{n-1} f\left(\frac{n\pi}{2Na}, \frac{m\pi}{2Na}\right) \delta_{p1} \delta_{r1} + \frac{1}{2} f\left(\frac{n\pi}{2Na}, \frac{n\pi}{2Na}\right) \delta_{p1} \delta_{r1} \right) \quad (144)$$

and in the case of the primitive-rectangular lattice

$$F_N = \frac{4A_{\text{BZ}}}{N^2} \sum_{n=0}^{2N} \sum_{m=0}^{2N} f\left(\frac{n\pi}{2Na}, \frac{m\pi}{2Nb}\right) \delta_{p1} \delta_{r1} \quad (145)$$

where $p \equiv n \pmod{2}$ and $r \equiv m \pmod{2}$. δ_{xy} is the Kronecker Delta where $\delta_{xy} = 1$ if $x = y$ and $\delta_{xy} = 0$ if $x \neq y$.

Here N is the number of points used along the $k_x(k_z)$ axis. In the irreducible segment of the square lattice, this is the largest number of points along that axis. In a primitive-rectangular lattice, no such distinction exists. a, b are lattice constants with $a = b$ in the case of the square lattice.

A.2 Author's 2D Brillouin Zone integration method

The method outlined in Subsection A.1 is the most efficient algorithm for integrating over the 2DBZ by minimising the number of regularly spaced points required. The form given by Cunningham requires the user to determine whether convergence has been achieved, repeating with more points if necessary. If the function is slow to converge, this can be quite time consuming. The author has made two developments to the results of Subsection A.1. In both cases 'automatic' techniques are used, where if convergence hasn't been met, new points are used without discarding the work already done. In each case, the algorithm is for the square lattice, though it is easily extendible (and easier to write) to the primitive-rectangular lattice.

The simplest of these two techniques, utilised in the header file `cunningham.h` is as follows, where in the case of the square lattice

$$F_{3^q N} = \frac{8A_{\text{BZ}}}{(3^q N)^2} \sum_{q=0}^{\alpha} \sum_{n=0}^{2(3^q)N} \left(\sum_{m=0}^{n-1} f\left(\frac{n\pi}{2(3^q)Na}, \frac{m\pi}{2(3^q)Na}\right) \delta_{p1} \delta_{r1} l + \frac{1}{2} f\left(\frac{n\pi}{2(3^q)Na}, \frac{n\pi}{2(3^q)Na}\right) \delta_{p1} \delta_{r1} l \right) \quad (146)$$

and in the case of the primitive-rectangular lattice

$$F_{3^q N} = \frac{4A_{\text{BZ}}}{(3^\alpha N)^2} \sum_{q=0}^{\alpha} \sum_{n=0}^{2(3^q)N} \sum_{m=0}^{2(3^q)N} f\left(\frac{n\pi}{2(3^q)Na}, \frac{m\pi}{2(3^q)Nb}\right) \delta_{p1} \delta_{r1} l \quad (147)$$

where $q = \alpha$ is an integer such that $\left|1 - \left|\frac{F_{2^{\alpha-1}N}}{F_{2^\alpha N}}\right|\right| \leq \delta$ and δ is the required relative error. $l = 1$ if $n \pmod{3} = 0 \wedge m \pmod{3} = 0$ and $l = 0$ otherwise. These equations are essentially revisions to equations (144) & (145). It can be observed that the combination of l and the summation over q means that points are 'recycled' rather than recalculated, hence optimising efficiency. N is the number of points used in each meridian at the start of the algorithm with different starting values leading to different completion times. Though per function, the optimal value N is only obtained by trial and error, it is recommended that the user starts with small N . The difference in completion time depending on N is not negligible; if ϕ is the time it takes to complete with $\alpha = c$ for some integer c , then 9ϕ is the total completion time when $\alpha = c+1$.

A case can be argued whether the efficiency of reusing the previously calculated points is favourable to simply quadrupling (in for example the primitive-rectangle) the number of points. By reusing the previous points, the next pass uses nine times as many points, as indicated above. Upon quadrupling the number of points and discarding the previous work, there is a possibility convergence may be achieved more quickly. However, this method penalises performance for allowing automatic convergence. Following the analysis of W. H. Press *et. al.* [57] and adapting to the square lattice indicates on average, the unnecessary work required for Equation (146) versus that of the quadrupling routine is a scaling of $\sqrt{\frac{15}{7}} \approx 1.46$, whereas the penalty incurred for not reusing the points in the quadrupling routine is variable, depending on the number of iterations and starting number of points, ranging from a performance factor of 1.66 to 1.33. The factor of 1.33 arises from a high number of starting points used, and one further pass being required. The factor of 1.66 arises from a low starting number of one point and a low number of passes required. For a low starting number and a high number of passes required, the factor converged to 1.5. Therefore the median ratio of 1.46 : 1.5 puts Equation (146) as more favourable. This is qualitatively useful, as it further emphasises the usefulness of being able to start with a small number of points, whereas in the alternative scheme, the penalty is compounded if multiple passes are required.

Applying the above analysis to the primitive-rectangle, the unnecessary work required for Equation (147) versus that of the quadrupling routine, is a scaling of 1.5 whereas the penalty incurred in the quadrupling method for not reusing points is a scaling of 1.33 on the number of evaluations required, if only one further pass is performed (no matter what the starting number of points is, in contrast with the square lattice). For a large number of passes, the total performance penalty converges to 1.5 thus in this case the total work required is on average equivalent. In this case, it is worth using the quadrupling technique if the integral is expected to converge quickly, and Equation (147) otherwise.

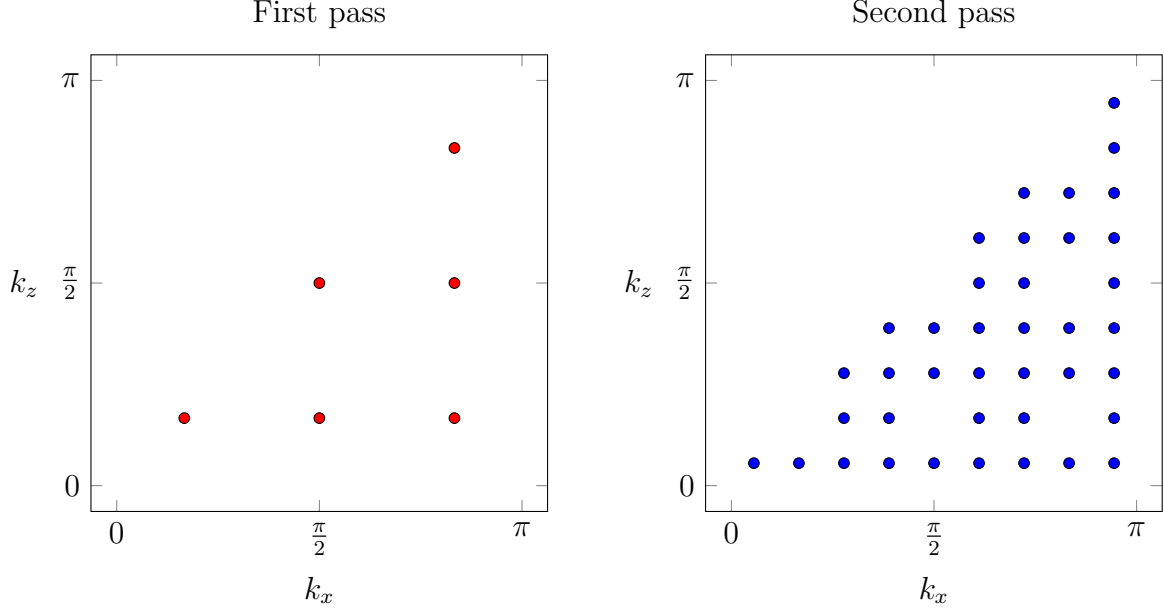


Figure 11: Author's 'automatic' Cunningham Special Points (CSP) routine with $N = 3$ and $a = 1$ showing the first pass on the left ($q = 0$) and the second pass on the right ($q = 1$). The figure illustrates how there is no penalty incurred when the number of points is increased, where the total number of points is the sum of both frames. Each point along the diagonal has half the weighting of the remaining points.

The coded routine requires the following input parameters;

- The starting number of points per meridian, N *must be an integer*
- The relative error tolerance, δ *in decimal form*
- The maximum number of points per meridian *must be an integer*
- The lattice constant, a
- Any number of further parameters to be forwarded to f

The second more sophisticated technique, utilised in the `cunningham_spawn.h` header file is an adaptive routine which, in the case of a primitive rectangular lattice, calls the integrand at the mean value point [52]. Then it calls the integrand at a further eight points, which together with the mean value point, forms a rectangular mesh with evenly spaced points such that if mirrored on the k_x or k_z axes, the mesh appears continuous. The routine then compares the initial point with the sum of the nine points scaled by $1/9$. If convergence is achieved, the function terminates. Otherwise, through recursion it samples each of these individual parent points and generates a sub-mesh of eight child points, with the parent point rescaled to form a ninth. The child points are positioned at a third of the distance to the next neighbour parent point, such that if each parent spawns children, again a continuous mesh is observed. The child points are summed and appropriately scaled to check for convergence against their corresponding parent. If convergence hasn't been met, the process repeats until convergence along all points is achieved, or the recursion limit has been reached.

It is easy to generalise this technique to the square lattice, where if a point is sampled on the line $k_x = k_z$ it is halved, and only five further child points are generated, since the further three do not lie within the irreducible segment. Two of these also lie on the line $k_x = k_z$, which means they are also weighted by a half. When spawning further points, only the parents that lie on the line $k_x = k_z$ spawn children that must also adhere to this special case. Otherwise the nine points as outlined in the previous paragraph are sampled. The routine requires the following input parameters;

- The recursion depth cap *must be an integer*
- The relative error tolerance *in decimal form*
- The lattice constant, a
- Any number of further parameters to be forwarded to f

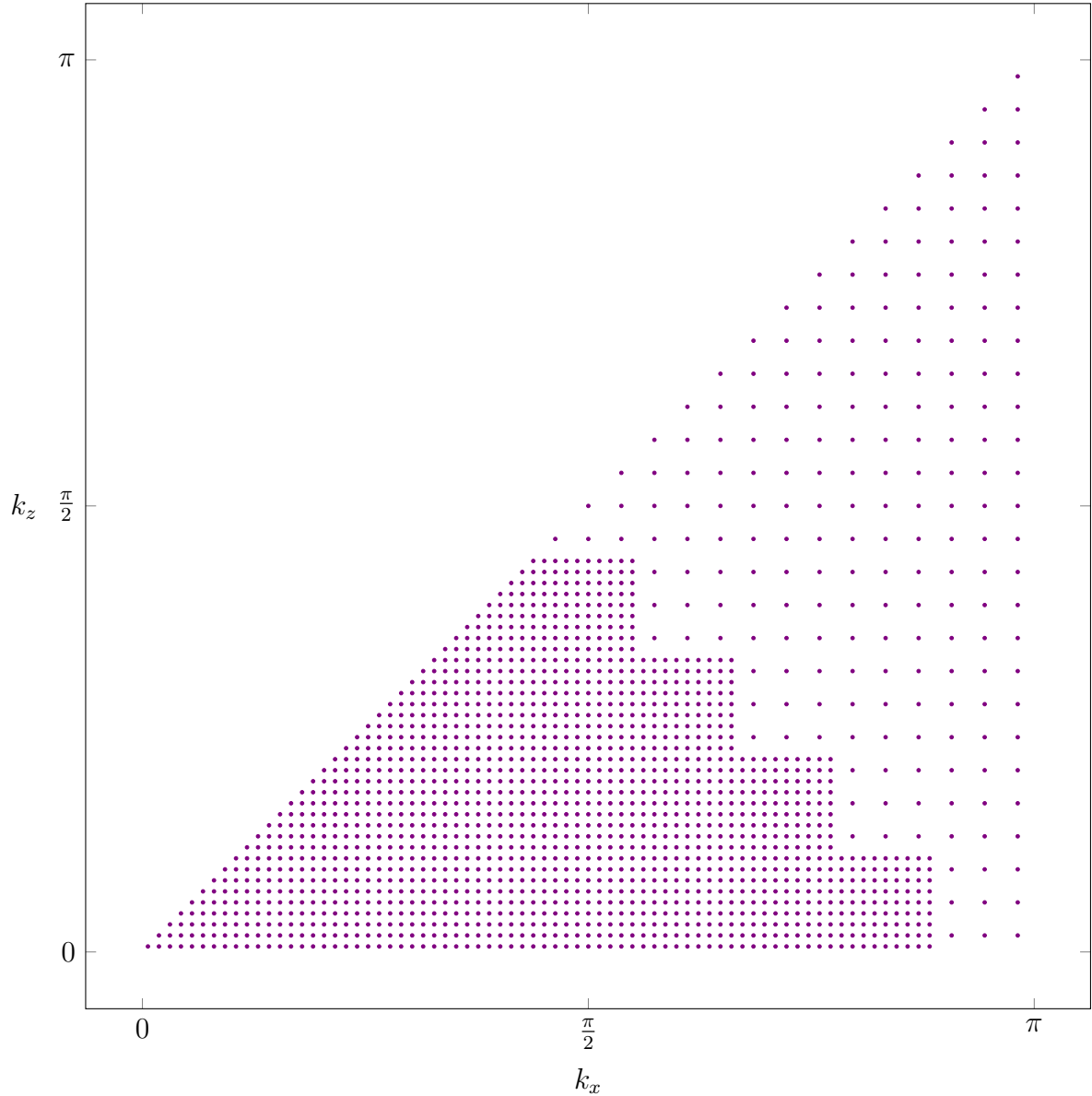


Figure 12: Typical sampling map of the author's automatic, adaptive CSP routine with $a = 1$ in the square lattice. Note the separation of points is smaller in some areas than others indicating convergence is slower in these places.

In each case, the routine has been concisely written as a header file within C++14 made callable for a number of different types including the Tux Family Eigen container type. The compiler is able to infer the type required in the routine from the return type of the function forwarded in the parameter list within the initial calling sequence. The routine makes use of perfect forwarding of r-value references to improve efficiency. Forwarding of an arbitrary number of parameters to f is made possible by means of variadic templates.

The following table shows benchmarking timings from the author's machine¹ be-

¹Intel Core i3-3120M CPU @ 2.50GHz, 4GB RAM

tween the first routine outlined in this section, `cunningham.h`, and the second routine `cunningham_spawn.h`. In each case, the SC one band model was calculated. Benchmarking against the original Cunningham Special Points (CSP) routine hasn't been attempted since it is difficult to create unbiased results as per the discussion over Equation (147) in this section.

Computation	Return type of integrand	Time to output results using <code>cunningham.h</code>	Time to output results using <code>cunningham_spawn.h</code>
Local density of states of a semi-infinite crystal (8 layers)	Real (double)	2.39 seconds	1.13 seconds
Exchange Coupling using summation of Matsubara Frequencies	Container of type complex double	1.60 seconds	0.91 seconds
Transport spincurrent using the Keldysh formalism, avoiding the energy integral	Container of type real (double)	81.39 seconds	34.08 seconds

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