PEETER JOOT
QUANTUM MECHANICSII

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 PEETER JOOTUofT PHY456H1F 2012 notes
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Peeter Joot: Quantum Mechanics II, UofT PHY456H1F 2012 notes, © August 2023

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I reserve the right to impose dictatorial control over any editing and content decisions, and may not accept merge requests as-is, or at all. That said, I will probably not refuse reasonable suggestions or merge requests.

## Dedicated to:

Aurora and Lance, my awesome kids, and
Sofia, who not only tolerates and encourages my studies, but is also awesome enough to think that math is sexy.

## PREFACE

These are my personal lecture notes for the Fall 2011, University of Toronto Quantum mechanics II course (PHY456H1F), taught by Prof. John E Sipe.

The official description of this course was:
Quantum dynamics in Heisenberg and Schrodinger Pictures; WKB approximation; Variational Method; Time-Independent Perturbation Theory; Spin; Addition of Angular Momentum; Time-Dependent Perturbation Theory; Scattering.

This document contains a few things

- My lecture notes.
- Notes from reading of the text [4]. This may include observations, notes on what seem like errors, and some solved problems.
- Different ways of tackling some of the assigned problems than the solution sets.
- Some personal notes exploring details that were not clear to me from the lectures.
- Some worked problems.

There were three main themes in this course, my notes for which can be found in

- Approximate methods and perturbation: chapter 1-chapter 6,
- Spin, angular momentum, and two particle systems: chapter 7-chapter 12,
- Scattering theory: chapter 13 -chapter 15 ,

Unlike some of my other course notes compilations, this one is short and contains few worked problems. It appears that I did most of my problem sets on paper and subsequently lost my solutions. There are also some major defects in these notes:

- There are plenty of places where things weren't clear, and there are still comments to followup on those issues to understand them.
- There is redundant content, from back to back lectures on materials that included review of the previous lecture notes.
- A lot of the stuff in the appendix (mostly personal notes and musings) should be merged into the appropriate lecture note chapters. Some work along those lines has been started, but that work was very preliminary.
- I reworked some ideas from the original lecture notes to make sense of them (in particular, adiabatic approximation theory), but then didn't go back and consolidate all the different notes for the topic into a single coherent unit.
- There were Mathematica notebooks for some of the topics with issues that I never did figure out.
- Lots of typos, bad spelling, and horrendous grammar.
- The indexing is very spotty.

Hopefully, despite these and other defects, these notes may be of some value to other students of Quantum Mechanics.

I'd like to thank Professor Sipe for teaching this course. I learned a lot and it provided a great foundation for additional study.

Peeter Joot peeterjoot@pm.me

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### 1.1 APPROXIMATE METHODS FOR FINDING ENERGY EIGENVALUES AND EIGENKETS.

In many situations one has a Hamiltonian $H$

$$
\begin{equation*}
H\left|\Psi_{n \alpha}\right\rangle=E_{n}\left|\Psi_{n \alpha}\right\rangle \tag{1.1}
\end{equation*}
$$

Here $\alpha$ is a "degeneracy index" (example: as in Hydrogen atom).

Why?

- Because this simplifies dynamics. Given an expansion of an initial state in some basis

$$
\begin{equation*}
|\Psi(0)\rangle=\sum_{n \alpha}\left|\Psi_{n \alpha}\right\rangle\left\langle\Psi_{n \alpha} \mid \Psi(0)\right\rangle=\sum_{n \alpha} c_{n \alpha}\left|\Psi_{n \alpha}\right\rangle, \tag{1.2}
\end{equation*}
$$

the time evolution of the state is

$$
\begin{align*}
|\Psi(t)\rangle & =e^{-i H t / \hbar}|\Psi(0)\rangle \\
& =\sum_{n \alpha} c_{n \alpha} e^{-i H t / \hbar}\left|\Psi_{n \alpha}\right\rangle  \tag{1.3}\\
& =\sum_{n \alpha} c_{n \alpha} e^{-i E_{n} t / \hbar}\left|\Psi_{n \alpha}\right\rangle .
\end{align*}
$$

- An "applied field", as depicted crudely in fig. 1.1, can often be thought of a driving the system from one eigenstate to another.
- In Statistical mechanics, at thermal equilibrium, the expectation value of an observable $O$ is

$$
\begin{equation*}
\langle O\rangle=\frac{\sum_{n \alpha}\left\langle\Psi_{n \alpha}\right| O\left|\Psi_{n \alpha}\right\rangle e^{-\beta E_{n}}}{Z} \tag{1.4}
\end{equation*}
$$

where

$$
\begin{equation*}
\beta=\frac{1}{k_{\mathrm{B}} T}, \tag{1.5}
\end{equation*}
$$

and

$$
\begin{equation*}
Z=\sum_{n \alpha} e^{-\beta E_{n}} . \tag{1.6}
\end{equation*}
$$



Figure 1.1: Stimulated emission.

### 1.2 VARIATIONAL PRINCIPLE.

Consider any ket

$$
\begin{equation*}
|\Psi\rangle=\sum_{n \alpha} c_{n \alpha}\left|\Psi_{n \alpha}\right\rangle \tag{1.7}
\end{equation*}
$$

perhaps not even normalized, where the coefficients

$$
\begin{equation*}
c_{n \alpha}=\left\langle\Psi_{n \alpha} \mid \Psi\right\rangle \tag{1.8}
\end{equation*}
$$

are unknown. We do know that

$$
\begin{equation*}
\langle\Psi \mid \Psi\rangle=\sum_{n \alpha}\left|c_{n \alpha}\right|^{2} \tag{1.9}
\end{equation*}
$$

and can relate the average energy to the ground state energy as follows

$$
\begin{align*}
\frac{\langle\Psi| H|\Psi\rangle}{\langle\Psi \mid \Psi\rangle} & =\frac{\sum_{n \alpha}\left|c_{n \alpha}\right|^{2} E_{n}}{\sum_{m \beta}\left|c_{m \beta}\right|^{2}} \\
& \geq \frac{\sum_{n \alpha}\left|c_{n \alpha}\right|^{2} E_{0}}{\sum_{m \beta}\left|c_{m \beta}\right|^{2}}  \tag{1.10}\\
& =E_{0}
\end{align*}
$$

So for any ket we can form the upper bound for the ground state energy

$$
\begin{equation*}
\frac{\langle\Psi| H|\Psi\rangle}{\langle\Psi \mid \Psi\rangle} \geq E_{0} \tag{1.11}
\end{equation*}
$$

There is a whole set of strategies based on estimating the ground state energy. This is called the Variational principle for ground state. See §24.2 in the text [4].

We define the functional

$$
\begin{equation*}
E[\Psi]=\frac{\langle\Psi| H|\Psi\rangle}{\langle\Psi \mid \Psi\rangle} \geq E_{0} . \tag{1.12}
\end{equation*}
$$

In particular, if $|\Psi\rangle=c\left|\Psi_{0}\right\rangle$ where $\left|\Psi_{0}\right\rangle$ is the normalized ground state, then

$$
\begin{equation*}
E\left[c \Psi_{0}\right]=E_{0} \tag{1.13}
\end{equation*}
$$

## Example 1.1: Hydrogen atom.

The solution of the Hydrogen problem is given by

$$
\begin{equation*}
\langle\mathbf{r}| H\left|\mathbf{r}^{\prime}\right\rangle=\mathcal{H} \delta^{3}\left(\mathbf{r}-\mathbf{r}^{\prime}\right) \tag{1.14}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathcal{H}=-\frac{\hbar^{2}}{2 \mu} \nabla^{2}-\frac{e^{2}}{r} \tag{1.15}
\end{equation*}
$$

and $\mu$ is the reduced mass.
We know the exact ground state energy

$$
\begin{equation*}
E_{0}=-R_{y} \tag{1.16}
\end{equation*}
$$

where

$$
\begin{equation*}
R_{y}=\frac{\mu e^{4}}{2 \hbar^{2}} \approx 13.6 \mathrm{eV} \tag{1.17}
\end{equation*}
$$

The ground state wavefunction, plotted in fig. 1.2, is

$$
\begin{equation*}
\left\langle\mathbf{r} \mid \Psi_{0}\right\rangle=\Phi_{100}(\mathbf{r})=\left(\frac{1}{\pi a_{0}^{3}}\right)^{1 / 2} e^{-r / a_{0}} \tag{1.18}
\end{equation*}
$$

where

$$
\begin{equation*}
a_{0}=\frac{\hbar^{2}}{\mu e^{2}} \approx 0.53 \AA \tag{1.19}
\end{equation*}
$$



Figure 1.2: Hydrogen ground state magnitude.
We can estimate by evaluating

$$
\begin{align*}
\langle\Psi| H|\Psi\rangle & =\int d^{3} \mathbf{r} \Psi^{*}(\mathbf{r})\left(-\frac{\hbar^{2}}{2 \mu} \nabla^{2}-\frac{e^{2}}{r}\right) \Psi(\mathbf{r})  \tag{1.20}\\
\langle\Psi \mid \Psi\rangle & =\int d^{3} \mathbf{r}|\Psi(\mathbf{r})|^{2}
\end{align*}
$$

Using the trial wave function $\Psi=e^{-\alpha r^{2}}$, plotted in fig. 1.3


Figure 1.3: Gaussian test function.
we want to calculate the ground state energy estimate associated with this $\alpha$ parameter

$$
\begin{equation*}
E[\Psi] \rightarrow E(\alpha) . \tag{1.21}
\end{equation*}
$$

Explicitly, this is

$$
\begin{equation*}
E(\alpha)=\frac{\int d^{3} \mathbf{r} e^{-\alpha r^{2}}\left(-\frac{\hbar^{2}}{2 \mu} \boldsymbol{\nabla}^{2}-\frac{e^{2}}{r}\right) e^{-\alpha r^{2}}}{\int d^{3} \mathbf{r} e^{-2 \alpha r^{2}}} \tag{1.22}
\end{equation*}
$$

It can apparently be shown that this results in

$$
\begin{equation*}
E(\alpha)=A \alpha-B \alpha^{1 / 2} \tag{1.23}
\end{equation*}
$$

where

$$
\begin{align*}
& A=\frac{3 \hbar^{2}}{2 \mu} \\
& B=2 e^{2}\left(\frac{2}{\pi}\right)^{1 / 2} \tag{1.24}
\end{align*}
$$

(I tried to calculate this, and get a different answer, so I've either taken bad notes, or calculated wrong, or what was on the board was not right.) Equation (1.23) is of the form fig. 1.4, with the minimum energy occuring at


Figure 1.4: Ground state energy for the Gaussian test function.

$$
\begin{equation*}
\alpha_{0}=\left(\frac{\mu e^{2}}{\hbar^{2}}\right) \frac{8}{9 \pi} \tag{1.25}
\end{equation*}
$$

So

$$
\begin{equation*}
E\left(\alpha_{0}\right)=-\frac{\mu e^{4}}{2 \hbar^{2}} \frac{8}{3 \pi}=-0.85 R_{y} \tag{1.26}
\end{equation*}
$$

maybe not too bad...

## Example 1.2: Helium atom.

Assume an infinite nuclear mass with nucleus charge $2 e$


Figure 1.5: Position vectors for the Helium electrons.
The problem that we want to solve is

$$
\begin{equation*}
\left(-\frac{\hbar^{2}}{2 m} \nabla_{1}^{2}-\frac{\hbar^{2}}{2 m} \nabla_{2}^{2}-\frac{2 e}{r}+\frac{e^{2}}{\left|\mathbf{r}_{1}-\mathbf{r}_{2}\right|}\right) \Psi_{0}\left(\mathbf{r}_{1}, \mathbf{r}_{2}\right)=E_{0} \Psi_{0}\left(\mathbf{r}_{1}, \mathbf{r}_{2}\right) \tag{1.27}
\end{equation*}
$$

where $\Psi_{0}\left(\mathbf{r}_{1}, \mathbf{r}_{2}\right)$ is the ground state wavefunction. Nobody can solve this problem. It is one of the simplest real problems in QM that cannot be solved exactly. Suppose that we neglected the electron, electron repulsion. Then

$$
\begin{equation*}
\Psi_{0}\left(\mathbf{r}_{1}, \mathbf{r}_{2}\right)=\bar{\Phi}_{100}\left(\mathbf{r}_{1}\right) \bar{\Phi}_{100}\left(\mathbf{r}_{2}\right) \tag{1.28}
\end{equation*}
$$

where

$$
\begin{equation*}
\left(-\frac{\hbar^{2}}{2 m} \nabla^{2}-\frac{2 e}{r}\right) \bar{\Phi}_{100}(\mathbf{r})=\epsilon \bar{\Phi}_{100}(\mathbf{r}) \tag{1.29}
\end{equation*}
$$

and where

$$
\begin{equation*}
\epsilon=-4 R_{y} \tag{1.30}
\end{equation*}
$$

and

$$
\begin{equation*}
R_{y}=\frac{m e^{4}}{2 \hbar^{2}} \tag{1.31}
\end{equation*}
$$

This is the solution to

$$
\begin{align*}
\left(-\frac{\hbar^{2}}{2 m} \nabla_{1}^{2}-\frac{\hbar^{2}}{2 m} \nabla_{2}^{2}-\frac{2 e}{r}\right) \Psi_{0}^{(0)}\left(\mathbf{r}_{1}, \mathbf{r}_{2}\right) & =E_{0} \Psi_{0}\left(\mathbf{r}_{1}, \mathbf{r}_{2}\right)  \tag{1.32}\\
& =E_{0}^{(0)} \Psi_{0}^{(0)}\left(\mathbf{r}_{1}, \mathbf{r}_{2}\right)
\end{align*}
$$

so

$$
\begin{equation*}
E_{0}^{(0)}=-8 R_{y} . \tag{1.33}
\end{equation*}
$$

Now we want to put back in the electron electron repulsion, and make an estimate.

With a trial wavefunction

$$
\begin{equation*}
\Psi\left(\mathbf{r}_{1}, \mathbf{r}_{2}, Z\right)=\left(\left(\frac{Z^{3}}{\pi a_{0}^{3}}\right)^{1 / 2} e^{-Z r_{1} / a_{0}}\right)\left(\left(\frac{Z^{3}}{\pi a_{0}^{3}}\right)^{1 / 2} e^{-Z r_{2} / a_{0}}\right) \tag{1.34}
\end{equation*}
$$

expect that the best estimate is for $Z \in[1,2]$. This can be calculated numerically, and we find

$$
\begin{equation*}
E(Z)=2 R_{Y}\left(Z^{2}-4 Z+\frac{5}{8} Z\right) \tag{1.35}
\end{equation*}
$$

The $Z^{2}$ comes from the kinetic energy. The $-4 Z$ is the electron nuclear attraction, and the final term is from the electron-electron repulsion.

The actual minimum is

$$
\begin{equation*}
Z=2-\frac{5}{16} \tag{1.36}
\end{equation*}
$$

$$
\begin{equation*}
E(2-5 / 16)=-77.5 \mathrm{eV} \tag{1.37}
\end{equation*}
$$

As expected, this is greater than the actual interaction free value of -78.6 eV .

## PERTURBATION METHODS.

### 2.1 States and wave functions.

Suppose we have the following non-degenerate energy eigenstates

$$
\begin{gather*}
\vdots \\
E_{1} \sim\left|\psi_{1}\right\rangle  \tag{2.1}\\
E_{0} \sim\left|\psi_{0}\right\rangle,
\end{gather*}
$$

and consider a state that is "very close" to $\left|\psi_{n}\right\rangle$

$$
\begin{equation*}
|\psi\rangle=\left|\psi_{n}\right\rangle+\left|\delta \psi_{n}\right\rangle . \tag{2.2}
\end{equation*}
$$

We form projections onto $\left|\psi_{n}\right\rangle$ "direction". The difference from this projection will be written $\left|\psi_{n \perp}\right\rangle$, as depicted in fig. 2.1. This illustration cannot not be interpreted literally, but illustrates the idea nicely. For the


Figure 2.1: Pictorial illustration of get projections.
amount along the projection onto $\left|\psi_{n}\right\rangle$ we write

$$
\begin{equation*}
\left\langle\psi_{n} \mid \delta \psi_{n}\right\rangle=\delta \alpha, \tag{2.3}
\end{equation*}
$$

so that the total deviation from the original state is

$$
\begin{equation*}
\left|\delta \psi_{n}\right\rangle=\delta \alpha\left|\psi_{n}\right\rangle+\left|\delta \psi_{n \perp}\right\rangle . \tag{2.4}
\end{equation*}
$$

The varied ket is then

$$
\begin{equation*}
|\psi\rangle=(1+\delta \alpha)\left|\psi_{n}\right\rangle+\left|\delta \psi_{n \perp}\right\rangle, \tag{2.5}
\end{equation*}
$$

where

$$
\begin{equation*}
(\delta \alpha)^{2},\left\langle\delta \psi_{n \perp} \mid \delta \psi_{n \perp}\right\rangle \ll 1 \tag{2.6}
\end{equation*}
$$

In terms of these projections our kets magnitude is

$$
\begin{align*}
\langle\psi \mid \psi\rangle= & \left(\left(1+\delta \alpha^{*}\right)\left\langle\psi_{n}\right|+\left\langle\delta \psi_{n \perp}\right|\right)\left((1+\delta \alpha)\left|\psi_{n}\right\rangle+\left|\delta \psi_{n \perp}\right\rangle\right) \\
= & |1+\delta \alpha|^{2}\left\langle\psi_{n} \mid \psi_{n}\right\rangle+\left\langle\delta \psi_{n \perp} \mid \delta \psi_{n \perp}\right\rangle  \tag{2.7}\\
& +\left(1+\delta \alpha^{*}\right)\left\langle\psi_{n} \mid \delta \psi_{n \perp}\right\rangle+(1+\delta \alpha)\left\langle\delta \psi_{n \perp} \mid \delta \psi_{n}\right\rangle .
\end{align*}
$$

Because $\left\langle\delta \psi_{n \perp} \mid \delta \psi_{n}\right\rangle=0$ this is

$$
\begin{equation*}
\langle\psi \mid \psi\rangle=|1+\delta \alpha|^{2}\left\langle\delta \psi_{n \perp} \mid \delta \psi_{n \perp}\right\rangle . \tag{2.8}
\end{equation*}
$$

Similarly for the energy expectation we have

$$
\begin{align*}
\langle\psi \mid \psi\rangle= & \left(\left(1+\delta \alpha^{*}\right)\left\langle\psi_{n}\right|+\left\langle\delta \psi_{n \perp}\right|\right) H\left((1+\delta \alpha)\left|\psi_{n}\right\rangle+\left|\delta \psi_{n \perp}\right\rangle\right) \\
= & |1+\delta \alpha|^{2} E_{n}\left\langle\psi_{n} \mid \psi_{n}\right\rangle+\left\langle\delta \psi_{n \perp} \mid H\right\rangle \delta \psi_{n \perp}  \tag{2.9}\\
& +\left(1+\delta \alpha^{*}\right) E_{n}\left\langle\psi_{n} \mid \delta \psi_{n \perp}\right\rangle+(1+\delta \alpha) E_{n}\left\langle\delta \psi_{n \perp} \mid \delta \psi_{n}\right\rangle,
\end{align*}
$$

or

$$
\begin{equation*}
\langle\psi| H|\psi\rangle=E_{n}|1+\delta \alpha|^{2}+\left\langle\delta \psi_{n \perp}\right| H\left|\delta \psi_{n \perp}\right\rangle . \tag{2.10}
\end{equation*}
$$

This gives

$$
\begin{align*}
E[\psi] & =\frac{\langle\psi| H|\psi\rangle}{\langle\psi \mid \psi\rangle} \\
& =\frac{E_{n}|1+\delta \alpha|^{2}+\left\langle\delta \psi_{n \perp}\right| H\left|\delta \psi_{n \perp}\right\rangle}{|1+\delta \alpha|^{2}\left\langle\delta \psi_{n \perp} \mid \delta \psi_{n \perp}\right\rangle} \\
& =\frac{E_{n}+\frac{\left\langle\delta \psi_{n}\right| H \mid\left\langle\delta \psi_{n \perp}\right\rangle}{|1+\delta \alpha|^{2}}}{1+\frac{\left\langle\delta \psi_{n \perp} \mid \delta \psi_{n}\right\rangle}{|1 \delta \alpha|^{2}}}  \tag{2.11}\\
& =E_{n}\left(1-\frac{\left\langle\delta \psi_{n \perp} \mid \delta \psi_{n \perp}\right\rangle}{|1+\delta \alpha|^{2}}+\cdots\right)+\cdots \\
& =E_{n}\left[1+O\left(\left(\delta \psi_{n \perp}\right)^{2}\right)\right],
\end{align*}
$$

where

$$
\begin{equation*}
\left(\delta \psi_{n \perp}\right)^{2} \sim\left\langle\delta \psi_{n \perp} \mid \delta \psi_{n \perp}\right\rangle \tag{2.12}
\end{equation*}
$$



Figure 2.2: Illustration of variation of energy with variation of Hamiltonian.

An interpretation is that "small errors" in $|\psi\rangle$ do not lead to large errors in $E[\psi]$. It is reasonably easy to get a good estimate and $E_{0}$, although it is reasonably hard to get a good estimate of $\left|\psi_{0}\right\rangle$. This is for the same reason, because $E[]$ is not terribly sensitive.

### 2.2 EXCited states.

$$
\begin{align*}
& E_{2} \sim\left|\psi_{2}\right\rangle  \tag{2.13}\\
& E_{1} \sim\left|\psi_{1}\right\rangle \\
& E_{0} \sim\left|\psi_{0}\right\rangle
\end{align*}
$$

Suppose we wanted an estimate of $E_{1}$. If we knew the ground state $\left|\psi_{0}\right\rangle$. For any trial $|\psi\rangle$ form

$$
\begin{equation*}
\left|\psi^{\prime}\right\rangle=|\psi\rangle-\left|\psi_{0}\right\rangle\left\langle\psi_{0} \mid \psi\right\rangle . \tag{2.14}
\end{equation*}
$$

We are taking out the projection of the ground state from an arbitrary trial function. For a state written in terms of the basis states, allowing for an $\alpha$ degeneracy

$$
\begin{equation*}
|\psi\rangle=c_{0}\left|\psi_{0}\right\rangle+\sum_{n>0, \alpha} c_{n \alpha}\left|\psi_{n \alpha}\right\rangle \tag{2.15}
\end{equation*}
$$

where

$$
\begin{equation*}
\left\langle\psi_{0} \mid \psi\right\rangle=c_{0} \tag{2.16}
\end{equation*}
$$

and

$$
\begin{equation*}
\left|\psi^{\prime}\right\rangle=\sum_{n>0, \alpha} c_{n \alpha}\left|\psi_{n \alpha}\right\rangle \tag{2.17}
\end{equation*}
$$

Note that there are some theorems that tell us that the ground state is generally non-degenerate. Expanding the energy gives

$$
\begin{align*}
E\left[\psi^{\prime}\right] & =\frac{\left\langle\psi^{\prime}\right| H\left|\psi^{\prime}\right\rangle}{\left\langle\psi^{\prime} \mid \psi^{\prime}\right\rangle} \\
& =\frac{\sum_{n>0, \alpha}\left|c_{n \alpha}\right|^{2} E_{n}}{\sum_{m>0, \beta}\left|c_{m \beta}\right|^{2}} \geq E_{1} \tag{2.18}
\end{align*}
$$

Often do not know the exact ground state, although we might have a guess $\left|\tilde{\psi}_{0}\right\rangle$.

Should we attempt a 2 nd order rejection of the form

$$
\begin{equation*}
\left|\psi^{\prime \prime}\right\rangle=|\psi\rangle-\left|\tilde{\psi}_{0}\right\rangle\left\langle\tilde{\psi}_{0} \mid \psi\right\rangle \tag{2.19}
\end{equation*}
$$

we cannot prove that

$$
\begin{equation*}
\frac{\left\langle\psi^{\prime \prime}\right| H\left|\psi^{\prime \prime}\right\rangle}{\left\langle\psi^{\prime \prime} \mid \psi^{\prime \prime}\right\rangle} \geq E_{1} \tag{2.20}
\end{equation*}
$$

But sometimes, even if you do not know the ground state $\left|\psi_{0}\right\rangle$, can choose trial kets $\left|\psi^{\prime \prime \prime}\right\rangle$ such that $\left\langle\psi_{0} \mid \psi^{\prime \prime \prime}\right\rangle=0$. In that case is it true that

$$
\begin{equation*}
\frac{\left\langle\psi^{\prime \prime \prime}\right| H\left|\psi^{\prime \prime \prime}\right\rangle}{\left\langle\psi^{\prime \prime \prime} \mid \psi^{\prime \prime \prime}\right\rangle} \geq E_{1} . \tag{2.21}
\end{equation*}
$$

Somewhat remarkably, this is often possible. We talked last time about the Hydrogen atom. In that case, you can guess that the excited state is in the $2 s$ orbital and is therefore orthogonal to the $1 s$ orbital.

### 2.3 PROBLEMS.

## Exercise 2.1 Harmonic oscillator. (2011 psl/pl)

Let $H_{o}$ indicate the Hamiltonian of a 1D harmonic oscillator with mass $m$ and frequency $\omega$

$$
\begin{equation*}
H_{o}=\frac{P^{2}}{2 m}+\frac{1}{2} m \omega^{2} X^{2} \tag{2.22}
\end{equation*}
$$

and denote the energy eigenstates by $|n\rangle$, where $n$ is the eigenvalue of the number operator.

1. Find $\langle n| X^{4}|n\rangle$.
2. Quadratic pertubation. Find the ground state energy of the Hamiltonian $H=H_{o}+\gamma X^{2}$. You may assume $\gamma>0$. [Hint: This is not a trick question.]
3. Linear pertubation. Find the ground state energy of the Hamiltonian $H=H_{o}-\alpha X$. [Hint: This is a bit harder than part 2 but not much. Try "completing the square."]

Answer for Exercise 2.1

Part 1. $X^{4} \quad$ Working through A we have now got enough context to attempt the first part of the question, calculation of

$$
\begin{equation*}
\langle n| X^{4}|n\rangle \tag{2.23}
\end{equation*}
$$

We have calculated things like this before, such as

$$
\begin{equation*}
\langle n| X^{2}|n\rangle=\frac{\hbar}{2 m \omega}\langle n|\left(a+a^{\dagger}\right)^{2}|n\rangle . \tag{2.24}
\end{equation*}
$$

To continue we need an exact relation between $|n\rangle$ and $|n \pm 1\rangle$. Recall that $a|n\rangle$ was an eigenstate of $a^{\dagger} a$ with eigenvalue $n-1$. This implies that the eigenstates $a|n\rangle$ and $|n-1\rangle$ are proportional

$$
\begin{equation*}
a|n\rangle=c_{n}|n-1\rangle \tag{2.25}
\end{equation*}
$$

or

$$
\begin{align*}
\langle n| a^{\dagger} a|n\rangle & =\left|c_{n}\right|^{2}\langle n-1 \mid n-1\rangle=\left|c_{n}\right|^{2} \\
& =n\langle n \mid n\rangle  \tag{2.26}\\
& =n,
\end{align*}
$$

so that

$$
\begin{equation*}
a|n\rangle=\sqrt{n}|n-1\rangle . \tag{2.27}
\end{equation*}
$$

Similarly let

$$
\begin{equation*}
a^{\dagger}|n\rangle=b_{n}|n+1\rangle \tag{2.28}
\end{equation*}
$$

or

$$
\begin{align*}
\langle n| a a^{\dagger}|n\rangle & =\left|b_{n}\right|^{2}\langle n-1 \mid n-1\rangle=\left|b_{n}\right|^{2} \\
\langle n|\left(1+a^{\dagger} a\right)|n\rangle & =  \tag{2.29}\\
& =1+n,
\end{align*}
$$

so that

$$
\begin{equation*}
a^{\dagger}|n\rangle=\sqrt{n+1}|n+1\rangle \tag{2.30}
\end{equation*}
$$

We can now return to eq. (2.23), and find

$$
\begin{equation*}
\langle n| X^{4}|n\rangle=\frac{\hbar^{2}}{4 m^{2} \omega^{2}}\langle n|\left(a+a^{\dagger}\right)^{4}|n\rangle . \tag{2.31}
\end{equation*}
$$

Consider half of this braket

$$
\begin{align*}
\left(a+a^{\dagger}\right)^{2}|n\rangle & =\left(a^{2}+\left(a^{\dagger}\right)^{2}+a^{\dagger} a+a a^{\dagger}\right)|n\rangle \\
& =\left(a^{2}+\left(a^{\dagger}\right)^{2}+a^{\dagger} a+\left(1+a^{\dagger} a\right)\right)|n\rangle \\
& =\left(a^{2}+\left(a^{\dagger}\right)^{2}+1+2 a^{\dagger} a\right)|n\rangle \\
& =\sqrt{n-1} \sqrt{n-2}|n-2\rangle+\sqrt{n+1} \sqrt{n+2}|n+2\rangle+|n\rangle+2 n|n\rangle \tag{2.32}
\end{align*}
$$

Squaring, utilizing the Hermitian nature of the $X$ operator

$$
\begin{align*}
\langle n| X^{4}|n\rangle & =\frac{\hbar^{2}}{4 m^{2} \omega^{2}}\left((n-1)(n-2)+(n+1)(n+2)+(1+2 n)^{2}\right) \\
& =\frac{\hbar^{2}}{4 m^{2} \omega^{2}}\left(6 n^{2}+4 n+5\right) \tag{2.33}
\end{align*}
$$

Part 2. Quadratic ground state Find the ground state energy of the Hamiltonian $H=H_{0}+\gamma X^{2}$ for $\gamma>0$. The new Hamiltonian has the form

$$
\begin{equation*}
H=\frac{P^{2}}{2 m}+\frac{1}{2} m\left(\omega^{2}+\frac{2 \gamma}{m}\right) X^{2}=\frac{P^{2}}{2 m}+\frac{1}{2} m \omega^{2} X^{2} \tag{2.34}
\end{equation*}
$$

where

$$
\begin{equation*}
\omega^{\prime}=\sqrt{\omega^{2}+\frac{2 \gamma}{m}} \tag{2.35}
\end{equation*}
$$

The energy states of the Hamiltonian are thus

$$
\begin{equation*}
E_{n}=\hbar \sqrt{\omega^{2}+\frac{2 \gamma}{m}}\left(n+\frac{1}{2}\right) \tag{2.36}
\end{equation*}
$$

and the ground state of the modified Hamiltonian $H$ is thus

$$
\begin{equation*}
E_{0}=\frac{\hbar}{2} \sqrt{\omega^{2}+\frac{2 \gamma}{m}} \tag{2.37}
\end{equation*}
$$

Part 3. Linear ground state Find the ground state energy of the Hamiltonian $H=H_{0}-\alpha X$. With a bit of play, this new Hamiltonian can be factored into

$$
\begin{equation*}
H=\hbar \omega\left(b^{\dagger} b+\frac{1}{2}\right)-\frac{\alpha^{2}}{2 m \omega^{2}}=\hbar \omega\left(b b^{\dagger}-\frac{1}{2}\right)-\frac{\alpha^{2}}{2 m \omega^{2}}, \tag{2.38}
\end{equation*}
$$

where

$$
\begin{align*}
b & =\sqrt{\frac{m \omega}{2 \hbar}} X+\frac{i P}{\sqrt{2 m \hbar \omega}}-\frac{\alpha}{\omega \sqrt{2 m \hbar \omega}}  \tag{2.39}\\
b^{\dagger} & =\sqrt{\frac{m \omega}{2 \hbar}} X-\frac{i P}{\sqrt{2 m \hbar \omega}}-\frac{\alpha}{\omega \sqrt{2 m \hbar \omega}} .
\end{align*}
$$

From eq. (2.38) we see that we have the same sort of commutator relationship as in the original Hamiltonian

$$
\begin{equation*}
\left[b, b^{\dagger}\right]=1, \tag{2.40}
\end{equation*}
$$

and because of this, all the preceding arguments follow unchanged with the exception that the energy eigenstates of this Hamiltonian are shifted by a constant

$$
\begin{equation*}
H|n\rangle=\left(\hbar \omega\left(n+\frac{1}{2}\right)-\frac{\alpha^{2}}{2 m \omega^{2}}\right)|n\rangle, \tag{2.41}
\end{equation*}
$$

where the $|n\rangle$ states are simultaneous eigenstates of the $b^{\dagger} b$ operator

$$
\begin{equation*}
b^{\dagger} b|n\rangle=n|n\rangle . \tag{2.42}
\end{equation*}
$$

The ground state energy is then

$$
\begin{equation*}
E_{0}=\frac{\hbar \omega}{2}-\frac{\alpha^{2}}{2 m \omega^{2}} . \tag{2.43}
\end{equation*}
$$

This makes sense. A translation of the entire position of the system should not effect the energy level distribution of the system, but we have set our reference potential differently, and have this constant energy adjustment to the entire system.

Exercise 2.2 Position expectation, spinless hydrogen. (2011 psl/p2)
Show that for all energy eigenstates $\left|\Phi_{n l m}\right\rangle$ of the (spinless) hydrogen atom, where as usual $n, l$, and $m$ are respectively the principal, azimuthal, and magnetic quantum numbers, we have

$$
\begin{equation*}
\left\langle\Phi_{n l m}\right| X\left|\Phi_{n l m}\right\rangle=\left\langle\Phi_{n l m}\right| Y\left|\Phi_{n l m}\right\rangle=\left\langle\Phi_{n l m}\right| Z\left|\Phi_{n l m}\right\rangle=0 . \tag{2.44}
\end{equation*}
$$

[Hint: Take note of the parity of the spherical harmonics (see "quick summary" notes on the spherical harmonics).]

Answer for Exercise 2.2
The summary sheet provides us with the wavefunction

$$
\begin{equation*}
\left\langle\mathbf{r} \mid \Phi_{n l m}\right\rangle=\frac{2}{n^{2} a_{0}^{3 / 2}} \sqrt{\frac{(n-l-1)!}{(n+l)!)^{3}}} F_{n l}\left(\frac{2 r}{n a_{0}}\right) Y_{l}^{m}(\theta, \phi), \tag{2.45}
\end{equation*}
$$

where $F_{n l}$ is a real valued function defined in terms of Lagueere polynomials. Working with the expectation of the $X$ operator to start with we have

$$
\begin{align*}
\left\langle\Phi_{n l m}\right| & X\left|\Phi_{n l m}\right\rangle \\
& =\int\left\langle\Phi_{n l m} \mid \mathbf{r}^{\prime}\right\rangle\left\langle\mathbf{r}^{\prime}\right| X|\mathbf{r}\rangle\left\langle\mathbf{r} \mid \Phi_{n l m}\right\rangle d^{3} \mathbf{r} d^{3} \mathbf{r}^{\prime} \\
& =\int\left\langle\Phi_{n l m} \mid \mathbf{r}^{\prime}\right\rangle \delta\left(\mathbf{r}-\mathbf{r}^{\prime}\right) r \sin \theta \cos \phi\left\langle\mathbf{r} \mid \Phi_{n l m}\right\rangle d^{3} \mathbf{r} d^{3} \mathbf{r}^{\prime} \\
& =\int \Phi_{n l m}^{*}(\mathbf{r}) r \sin \theta \cos \phi \Phi_{n l m}(\mathbf{r}) d^{3} \mathbf{r} \\
& \sim \int r^{2} d r\left|F_{n l}\left(\frac{2 r}{n a_{0}}\right)\right|^{2} r \int \sin \theta d \theta d \phi Y_{l}^{m *}(\theta, \phi) \sin \theta \cos \phi Y_{l}^{m}(\theta, \phi) \tag{2.46}
\end{align*}
$$

Recalling that the only $\phi$ dependence in $Y_{l}^{m}$ is $e^{i m \phi}$ we can perform the $d \phi$ integration directly, which is

$$
\begin{equation*}
\int_{\phi=0}^{2 \pi} \cos \phi d \phi e^{-i m \phi} e^{i m \phi}=0 \tag{2.47}
\end{equation*}
$$

We have the same story for the $Y$ expectation which is

$$
\begin{align*}
& \left\langle\Phi_{n l m}\right| X\left|\Phi_{n l m}\right\rangle \\
&  \tag{2.48}\\
& \quad \sim \int r^{2} d r\left|F_{n l}\left(\frac{2 r}{n a_{0}}\right)\right|^{2} r \times \int \sin \theta d \theta d \phi Y_{l}^{m *}(\theta, \phi) \sin \theta \sin \phi Y_{l}^{m}(\theta, \phi)
\end{align*}
$$

Our $\phi$ integral is then just

$$
\begin{equation*}
\int_{\phi=0}^{2 \pi} \sin \phi d \phi e^{-i m \phi} e^{i m \phi}=0 \tag{2.49}
\end{equation*}
$$

also zero. The $Z$ expectation is a slightly different story. There we have

$$
\begin{align*}
\left\langle\Phi_{n l m}\right| Z \mid & \left|\Phi_{n l m}\right\rangle \\
\sim & \int d r\left|F_{n l}\left(\frac{2 r}{n a_{0}}\right)\right|^{2} r^{3} \\
& \int_{0}^{2 \pi} d \phi \int_{0}^{\pi} \sin \theta d \theta(\sin \theta)^{-2 m}\left(\frac{d^{l-m}}{d(\cos \theta)^{l-m}} \sin ^{2 l} \theta\right)^{2} \cos \theta \tag{2.50}
\end{align*}
$$

Within this last integral we can make the substitution

$$
\begin{align*}
u & =\cos \theta \\
\sin \theta d \theta & =-d(\cos \theta)=-d u  \tag{2.51}\\
u & \in[1,-1]
\end{align*}
$$

and the integral takes the form

$$
\begin{equation*}
-\int_{-1}^{1}(-d u) \frac{1}{\left(1-u^{2}\right)^{m}}\left(\frac{d^{l-m}}{d u^{l-m}}\left(1-u^{2}\right)^{l}\right)^{2} u \tag{2.52}
\end{equation*}
$$

Here we have the product of two even functions, times one odd function ( $u$ ), over a symmetric interval, so the end result is zero, completing the problem. I was not able to see how to exploit the parity result suggested in the problem, but it was not so bad to show these directly.

Exercise 2.3 Angular momentum operator. (2011 psl/p3)
Working with the appropriate expressions in Cartesian components, confirm that $L_{i}|\psi\rangle=0$ for each component of angular momentum $L_{i}$, if $\langle\mathbf{r} \mid \psi\rangle=\psi(\mathbf{r})$ is in fact only a function of $r=|\mathbf{r}|$.
Answer for Exercise 2.3
In order to proceed, we will have to consider a matrix element, so that we can operate on $|\psi\rangle$ in position space. For that matrix element, we can
proceed to insert complete states, and reduce the problem to a question of wavefunctions. That is

$$
\begin{align*}
\langle\mathbf{r}| L_{i}|\psi\rangle & =\int d^{3} \mathbf{r}^{\prime}\langle\mathbf{r}| L_{i}\left|\mathbf{r}^{\prime}\right\rangle\left\langle\mathbf{r}^{\prime} \mid \psi\right\rangle \\
& =\int d^{3} \mathbf{r}^{\prime}\langle\mathbf{r}| \epsilon_{i a b} X_{a} P_{b}\left|\mathbf{r}^{\prime}\right\rangle\left\langle\mathbf{r}^{\prime} \mid \psi\right\rangle \\
& =-i \hbar \epsilon_{i a b} \int d^{3} \mathbf{r}^{\prime} x_{a}\langle\mathbf{r}| \frac{\partial \psi\left(\mathbf{r}^{\prime}\right)}{\partial X_{b}}\left|\mathbf{r}^{\prime}\right\rangle \\
& =-i \hbar \epsilon_{i a b} \int d^{3} \mathbf{r}^{\prime} x_{a} \frac{\partial \psi\left(\mathbf{r}^{\prime}\right)}{\partial x_{b}}\left\langle\mathbf{r} \mid \mathbf{r}^{\prime}\right\rangle  \tag{2.53}\\
& =-i \hbar \epsilon_{i a b} \int d^{3} \mathbf{r}^{\prime} x_{a} \frac{\partial \psi\left(\mathbf{r}^{\prime}\right)}{\partial x_{b}} \delta^{3}\left(\mathbf{r}-\mathbf{r}^{\prime}\right) \\
& =-i \hbar \epsilon_{i a b} x_{a} \frac{\partial \psi(\mathbf{r})}{\partial x_{b}}
\end{align*}
$$

With $\psi(\mathbf{r})=\psi(r)$ we have

$$
\begin{align*}
\langle\mathbf{r}| L_{i}|\psi\rangle & =-i \hbar \epsilon_{i a b} x_{a} \frac{\partial \psi(r)}{\partial x_{b}} \\
& =-i \hbar \epsilon_{i a b} x_{a} \frac{\partial r}{\partial x_{b}} \frac{d \psi(r)}{d r}  \tag{2.54}\\
& =-i \hbar \epsilon_{i a b} x_{a} \frac{1}{2} 2 x_{b} \frac{1}{r} \frac{d \psi(r)}{d r} .
\end{align*}
$$

We are left with an sum of a symmetric product $x_{a} x_{b}$ with the antisymmetric tensor $\epsilon_{i a b}$ so this is zero for all $i \in[1,3]$.

## Exercise 2.4 Helium atom, variational method, first steps.

Verify (24.69) from [4], calculating the required six fold integral. Answer for Exercise 2.4

To verify (24.69) from [4], a six fold integral is required

$$
\begin{align*}
\left\langle-\frac{\hbar^{2}}{2 m}\left(\nabla_{1}^{2}+\nabla_{2}^{2}\right)\right\rangle= & -\frac{\hbar^{2}}{2 m} \int d r_{1} d \Omega_{1} r_{1}^{2} d r_{2} d \Omega_{2} r_{2}^{2} \frac{Z^{6}}{\pi^{2} a_{0}^{6}} e^{-\left(r_{1}+r_{2}\right) Z / a_{0}} \\
& \left(\frac{2}{r_{1}} \frac{\partial}{\partial r_{1}}+\frac{\partial^{2}}{\partial r_{1}^{2}}+\frac{2}{r_{2}} \frac{\partial}{\partial r_{2}}+\frac{\partial^{2}}{\partial r_{2}^{2}}\right) e^{-\left(r_{1}+r_{2}\right) Z / a_{0}} \\
= & -\frac{\hbar^{2}}{2 m} \frac{Z^{6}}{\pi^{2} a_{0}^{6}}(4 \pi)^{2} \int d r_{1} d r_{2} r_{1}^{2} r_{2}^{2} e^{-\left(r_{1}+r_{2}\right) Z / a_{0}}  \tag{2.55}\\
& \left(\frac{2}{r_{1}} \frac{\partial}{\partial r_{1}}+\frac{\partial^{2}}{\partial r_{1}^{2}}+\frac{2}{r_{2}} \frac{\partial}{\partial r_{2}}+\frac{\partial^{2}}{\partial r_{2}^{2}}\right) e^{-\left(r_{1}+r_{2}\right) Z / a_{0}}
\end{align*}
$$

Making a change of variables

$$
\begin{align*}
& x=\frac{Z r_{1}}{a_{0}} \\
& y=\frac{Z r_{2}}{a_{0}} \tag{2.56}
\end{align*}
$$

we have

$$
\begin{align*}
\left\langle-\frac{\hbar^{2}}{2 m}\right. & \left.\left(\nabla_{1}^{2}+\nabla_{2}^{2}\right)\right\rangle \\
& =-\frac{8 \hbar^{2}}{m} \frac{Z^{2}}{a_{0}^{2}} \int d x d y x^{2} y^{2} e^{-x-y}\left(\frac{2}{x} \frac{\partial}{\partial x}+\frac{\partial^{2}}{\partial x^{2}}+\frac{2}{y} \frac{\partial}{\partial y}+\frac{\partial^{2}}{\partial y^{2}}\right) e^{-x-y} \\
& =-\frac{8 \hbar^{2}}{m} \frac{Z^{2}}{a_{0}^{2}} \int d x d y x^{2} y^{2} e^{-x-y}\left(-\frac{2}{x}+1-\frac{2}{y}+1\right) e^{-x-y} \\
& =\frac{16 \hbar^{2}}{m} \frac{Z^{2}}{a_{0}^{2}} \int d x d y x^{2} y^{2} e^{-2 x-2 y}\left(\frac{1}{x}+\frac{1}{y}-1\right) \\
& =\frac{\hbar^{2}}{m} \frac{Z^{2}}{a_{0}^{2}} \tag{2.57}
\end{align*}
$$

With

$$
\begin{equation*}
a_{0}=\frac{\hbar^{2}}{m e^{2}} \tag{2.58}
\end{equation*}
$$

we have the result from the text

$$
\begin{equation*}
\left\langle-\frac{\hbar^{2}}{2 m}\left(\boldsymbol{\nabla}_{1}^{2}+\nabla_{2}^{2}\right)\right\rangle=\frac{Z^{2} e^{2}}{a_{0}} \tag{2.59}
\end{equation*}
$$

Verification of (24.70) follows in a similar fashion. We have

$$
\begin{align*}
& \left\langle 2 e^{2}\left(\frac{1}{r_{1}}+\frac{1}{r_{2}}\right)\right\rangle \\
& \quad=2 e^{2} \frac{Z^{6}}{\pi^{2} a_{0}^{6}}(4 \pi)^{2} \int e^{-2\left(r_{1}+r_{2}\right) Z / a_{0}} r_{1}^{2} r_{2}^{2} d r_{1} d r_{2}\left(\frac{1}{r_{1}}+\frac{1}{r_{2}}\right)  \tag{2.60}\\
& \quad=32 e^{2} \frac{Z}{a_{0}} \int e^{-2 x-2 y} x^{2} y^{2} d x d y\left(\frac{1}{x}+\frac{1}{y}\right) \\
& \quad=4 e^{2} \frac{Z}{a_{0}}
\end{align*}
$$

## Exercise 2.5 Helium expectation.

In §24.2.1 of the text [4] is an expectation value calculation associated with the Helium atom. Show that equation (24.76) is wrong.
Answer for Exercise 2.5
The following hand calculation and a Mathematica calculation shows that §24.2.1 [4] is wrong, provided no compensating error was made by me. Start with

$$
\begin{align*}
\left\langle\frac{e^{2}}{\left|\mathbf{r}_{1}-\mathbf{r}_{2}\right|}\right\rangle= & \left(\frac{Z^{3}}{\pi a_{0}^{3}}\right)^{2} e^{2} \int d^{3} k d^{3} r_{1} d^{3} r_{2} \frac{1}{2 \pi^{2} k^{2}} e^{i \mathbf{k} \cdot\left(\mathbf{r}_{1}-\mathbf{r}_{2}\right)} e^{-2 Z\left(r_{1}+r_{2}\right) / a_{0}} \\
= & \left(\frac{Z^{3}}{\pi a_{0}^{3}}\right)^{2} e^{2} \frac{1}{2 \pi^{2}} \int d^{3} k \frac{1}{k^{2}} \int d^{3} r_{1} e^{i \mathbf{k} \cdot \mathbf{r}_{1}} e^{-2 Z r_{1} / a_{0}} \times \\
& \int d^{3} r_{2} e^{-i \mathbf{k} \cdot \mathbf{r}_{2}} e^{-2 Z r_{2} / a_{0}} \tag{2.61}
\end{align*}
$$

To evaluate the two last integrals, I figure the author has aligned the axis for the $d^{3} r_{1}$ volume elements to make the integrals easier. Specifically, for the first so that $\mathbf{k} \cdot \mathbf{r}_{1}=k r_{1} \cos \theta$, so the integral takes the form

$$
\begin{align*}
\int d^{3} r_{1} e^{i \mathbf{k} \cdot \mathbf{r}_{1}} e^{-2 Z r_{1} / a_{0}} & =-\int r_{1}^{2} d r_{1} d \phi d(\cos \theta) e^{i k r_{1} \cos \theta} e^{-2 Z r_{1} / a_{0}} \\
& =-2 \pi \int_{r=0}^{\infty} \int_{u=1}^{-1} r^{2} d r d u e^{i k r u} e^{-2 Z r / a_{0}} \\
& =-2 \pi \int_{r=0}^{\infty} r^{2} d r \frac{1}{i k r}\left(e^{-i k r}-e^{i k r}\right) e^{-2 Z r / a_{0}}  \tag{2.62}\\
& =\frac{4 \pi}{k} \int_{r=0}^{\infty} r d r \frac{1}{2 i}\left(e^{i k r}-e^{-i k r}\right) e^{-2 Z r / a_{0}} \\
& =\frac{4 \pi}{k} \int_{r=0}^{\infty} r d r \sin (k r) e^{-2 Z r / a_{0}}
\end{align*}
$$

For this last, Mathematica gives me (24.75) from the text

$$
\begin{equation*}
\int d^{3} r_{1} e^{i \mathbf{k} \cdot \mathbf{r}_{1}} e^{-2 Z r_{1} / a_{0}}=\frac{16 \pi Z a_{0}^{3}}{\left(k^{2} a_{0}^{2}+4 Z^{2}\right)^{2}} \tag{2.63}
\end{equation*}
$$

For the second integral, if we align the axis so that $-\mathbf{k} \cdot \mathbf{r}_{2}=k r \cos \theta$ and repeat, then we have

$$
\begin{align*}
\left\langle\frac{e^{2}}{\left|\mathbf{r}_{1}-\mathbf{r}_{2}\right|}\right\rangle & =\left(\frac{Z^{3}}{\pi a_{0}^{3}}\right)^{2} e^{2} \frac{1}{2 \pi^{2}} 16^{2} \pi^{2} Z^{2} a_{0}^{6} \int d^{3} k \frac{1}{k^{2}} \frac{1}{\left(k^{2} a_{0}^{2}+4 Z^{2}\right)^{4}} \\
& =\frac{128 Z^{8}}{\pi^{2}} e^{2} \int d k d \Omega \frac{1}{\left(k^{2} a_{0}^{2}+4 Z^{2}\right)^{4}}  \tag{2.64}\\
& =\frac{512 Z^{8}}{\pi} e^{2} \int d k \frac{1}{\left(k^{2} a_{0}^{2}+4 Z^{2}\right)^{4}} \\
& =\frac{512 Z^{8}}{\pi} e^{2} \int d k \frac{1}{\left(k^{2} a_{0}^{2}+4 Z^{2}\right)^{4}}
\end{align*}
$$

With $k a_{0}=2 Z_{k}$ this is

$$
\begin{align*}
\left\langle\frac{e^{2}}{\left|\mathbf{r}_{1}-\mathbf{r}_{2}\right|}\right\rangle & =\frac{512 Z^{8}}{\pi} e^{2} \int d \kappa \frac{2 Z}{a_{0}} \frac{1}{(2 Z)^{8}\left(\kappa^{2}+1\right)^{4}} \\
& =\frac{4 Z}{\pi a_{0}} e^{2} \int d \kappa \frac{1}{\left(\kappa^{2}+1\right)^{4}}  \tag{2.65}\\
& =\frac{4 Z}{\pi a_{0}} e^{2} \frac{5 \pi}{32}
\end{align*}
$$

This gives us

$$
\begin{equation*}
\left\langle\frac{e^{2}}{\left|\mathbf{r}_{1}-\mathbf{r}_{2}\right|}\right\rangle=\frac{4 Z e^{2}}{\pi a_{0}} \frac{5 \pi}{32}=\frac{5 Z e^{2}}{8 a_{0}} \tag{2.66}
\end{equation*}
$$

## Exercise 2.6 SHO ground state energy, variational method.

Problem 3 of $\S 24.4$ in the text [4] asks to use the variational method to find the ground state energy of a one dimensional harmonic oscillator Hamiltonian.
Answer for Exercise 2.6

Recap. Variational method to find the ground state energy. Somewhat unexpectedly, once I take derivatives equate to zero, I find that the variational parameter beta becomes imaginary?

I tried this twice on paper and pencil, both times getting the same thing. This seems like a noteworthy problem, and one worth reflecting on a bit.

Recap. The variational method. Given any, not necessarily normalized wavefunction, with a series representation specified using the energy eigenvectors for the space

$$
\begin{equation*}
|\psi\rangle=\sum_{m} c_{m}\left|\psi_{m}\right\rangle, \tag{2.67}
\end{equation*}
$$

where

$$
\begin{equation*}
H\left|\psi_{m}\right\rangle=E_{m}\left|\psi_{m}\right\rangle, \tag{2.68}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\langle\psi_{m} \mid \psi_{n}\right\rangle=\delta_{m n} . \tag{2.69}
\end{equation*}
$$

We can perform an energy expectation calculation with respect to this more general state

$$
\begin{align*}
\langle\psi| H|\psi\rangle & =\sum_{m} c_{m}^{*}\left\langle\psi_{m}\right| H \sum_{n} c_{n}\left|\psi_{n}\right\rangle \\
& =\sum_{m} c_{m}^{*}\left\langle\psi_{m}\right| \sum_{n} c_{n} E_{n}\left|\psi_{n}\right\rangle \\
& =\sum_{m, n} c_{m}^{*} c_{n} E_{n}\left\langle\psi_{m} \mid \psi_{n}\right\rangle  \tag{2.70}\\
& \sum_{m}\left|c_{m}\right|^{2} E_{m} \\
& \geq \sum_{m}\left|c_{m}\right|^{2} E_{0} \\
& =E_{0}\langle\psi \mid \psi\rangle
\end{align*}
$$

This allows us to form an estimate of the ground state energy for the system, by using any state vector formed from a superposition of energy eigenstates, by simply calculating

$$
\begin{equation*}
E_{0} \leq \frac{\langle\psi| H|\psi\rangle}{\langle\psi \mid \psi\rangle} . \tag{2.71}
\end{equation*}
$$

One of the examples in the text is to use this to find an approximation of the ground state energy for the Helium atom Hamiltonian

$$
\begin{equation*}
H=-\frac{\hbar^{2}}{2 m}\left(\boldsymbol{\nabla}_{1}^{2}+\boldsymbol{\nabla}_{1}^{2}\right)-2 e^{2}\left(\frac{1}{r_{1}}+\frac{1}{r_{2}}\right)+\frac{e^{2}}{\left|\mathbf{r}_{1}-\mathbf{r}_{2}\right|} . \tag{2.72}
\end{equation*}
$$

This calculation is performed using a trial function that was a solution of the interaction free Hamiltonian

$$
\begin{equation*}
\phi=\frac{Z^{3}}{\pi a_{0}^{3}} e^{-Z\left(r_{1}+r_{2}\right) / a_{0}} . \tag{2.73}
\end{equation*}
$$

This is despite the fact that this is not a solution to the interaction Hamiltonian. The end result ends up being pretty close to the measured value (although there is a pesky error in the book that appears to require a compensating error somewhere else).

Part of the variational technique used in that problem, is to allow Z to vary, and then once the normalized expectation is computed, set the derivative of that equal to zero to calculate the trial wavefunction as a parameter of $Z$ that has the lowest energy eigenstate for a function of that form. We find considering the Harmonic oscillator that this final variation does not necessarily produce meaningful results.

The Harmonic oscillator variational problem. The problem asks for the use of the trial wavefunction

$$
\begin{equation*}
\phi=e^{-\beta|x|} \tag{2.74}
\end{equation*}
$$

to perform the variational calculation above for the Harmonic oscillator Hamiltonian, which has the one dimensional position space representation

$$
\begin{equation*}
H=-\frac{\hbar^{2}}{2 m} \frac{d^{2}}{d x^{2}}+\frac{1}{2} m \omega^{2} x^{2} \tag{2.75}
\end{equation*}
$$

We can find the normalization easily

$$
\begin{align*}
\langle\phi \mid \phi\rangle & =\int_{-\infty}^{\infty} e^{-2 \beta|x|} d x \\
& =2 \frac{1}{2 \beta} \int_{0}^{\infty} e^{-2 \beta x} 2 \beta d x  \tag{2.76}\\
& =2 \frac{1}{2 \beta} \int_{0}^{\infty} e^{-u} d u \\
& =\frac{1}{\beta} .
\end{align*}
$$

Using integration by parts, we find for the energy expectation
$\langle\phi| H|\phi\rangle$

$$
\begin{align*}
& =\int_{-\infty}^{\infty} d x e^{-\beta|x|}\left(-\frac{\hbar^{2}}{2 m} \frac{d^{2}}{d x^{2}}+\frac{1}{2} m \omega^{2} x^{2}\right) e^{-\beta|x|} \\
& =\lim _{\epsilon \rightarrow 0}\left(\int_{-\infty}^{-\epsilon}+\int_{-\epsilon}^{\epsilon}+\int_{\epsilon}^{\infty}\right) d x e^{-\beta|x|}\left(-\frac{\hbar^{2}}{2 m} \frac{d^{2}}{d x^{2}}+\frac{1}{2} m \omega^{2} x^{2}\right) e^{-\beta|x|} \\
& =2 \int_{0}^{\infty} d x e^{-2 \beta x}\left(-\frac{\hbar^{2} \beta^{2}}{2 m}+\frac{1}{2} m \omega^{2} x^{2}\right)-\frac{\hbar^{2}}{2 m} \lim _{\epsilon \rightarrow 0} \int_{-\epsilon}^{\epsilon} d x e^{-\beta|x|} \frac{d^{2}}{d x^{2}} e^{-\beta|x|} \tag{2.77}
\end{align*}
$$

The first integral we can do

$$
\begin{align*}
2 \int_{0}^{\infty} d x e^{-2 \beta x} & \left(-\frac{\hbar^{2} \beta^{2}}{2 m}+\frac{1}{2} m \omega^{2} x^{2}\right) \\
& =-\frac{\hbar^{2} \beta^{2}}{m} \int_{0}^{\infty} d x e^{-2 \beta x}+m \omega^{2} \int_{0}^{\infty} d x x^{2} e^{-2 \beta x} \\
& =-\frac{\hbar^{2} \beta}{2 m} \int_{0}^{\infty} d u e^{-u}+\frac{m \omega^{2}}{8 \beta^{3}} \int_{0}^{\infty} d u u^{2} e^{-u}  \tag{2.78}\\
& =-\frac{\beta \hbar^{2}}{2 m}+\frac{m \omega^{2}}{4 \beta^{3}}
\end{align*}
$$

A naive evaluation of this integral requires the origin to be avoided where the derivative of $|x|$ becomes undefined. This also provides a nice way to evaluate this integral because we can double the integral and half the range, eliminating the absolute value.

However, can we assume that the remaining integral is zero?
I thought that we could, but the end result is curious. I also verified my calculation symbolically in 24.4.3_attempt_with_mathematica.nb, but found that Mathematica required some special hand holding to deal with the origin. Initially I coded this by avoiding the origin as above, but later switched to $|x|=\sqrt{x^{2}}$ which Mathematica treats more gracefully.

Without that last integral, involving our singular $|x|^{\prime}$ and $|x|^{\prime \prime}$ terms, our ground state energy estimation, parameterized by $\beta$ is

$$
\begin{equation*}
E[\beta]=-\frac{\beta^{2} \hbar^{2}}{2 m}+\frac{m \omega^{2}}{4 \beta^{2}} . \tag{2.79}
\end{equation*}
$$

Observe that if we set the derivative of this equal to zero to find the "best" beta associated with this trial function

$$
\begin{equation*}
0=\frac{\partial E}{\partial \beta}=-\frac{\beta \hbar^{2}}{2 m}-\frac{m \omega^{2}}{2 \beta^{3}}, \tag{2.80}
\end{equation*}
$$

we find that the parameter beta that best minimizes this ground state energy function is complex with value

$$
\begin{equation*}
\beta^{2}= \pm \frac{i m \omega}{\sqrt{2} \hbar} . \tag{2.81}
\end{equation*}
$$

It appears at first glance that we can not minimize eq. (2.79) to find a best ground state energy estimate associated with the trial function eq. (2.74). We do however, know the exact ground state energy $\hbar \omega / 2$ for the Harmonic oscillator. It is possible to show that for all $\beta^{2}$ we have

$$
\begin{equation*}
\frac{\hbar \omega}{2} \leq-\frac{\beta^{2} \hbar^{2}}{2 m}+\frac{m \omega^{2}}{4 \beta^{2}} ? \tag{2.82}
\end{equation*}
$$

This inequality would be expected if we can assume that the trial wavefunction has a Fourier series representation utilizing the actual energy eigenfunctions for the system.

The resolution to this question is avoided once we include the singularity. This is explored in the last part of these notes.

Is our trial function representable? I thought perhaps that since the trial wave function for this problem lies outside the span of the Hilbert space that describes the solutions to the Harmonic oscillator. Another thing of possible interest is the trouble near the origin for this wave function, when operated on by $P^{2} / 2 m$, and this has been (incorrectly assumed to have zero contribution above).

I had initially thought that part of the value of this variational method was that we can use it despite not even knowing what the exact solution is (and in the case of the Helium atom, I believe it was stated in class that an exact closed form solution is not even known). This makes me wonder what restrictions must be imposed on the trial solutions to get a meaningful answer from the variational calculation?

Suppose that the trial wavefunction is not representable in the solution space. If that is the case, we need to adjust the treatment to account for that. Suppose we have

$$
\begin{equation*}
|\phi\rangle=\sum_{n} c_{n}\left|\psi_{n}\right\rangle+c_{\perp}\left|\psi_{\perp}\right\rangle \tag{2.83}
\end{equation*}
$$

where $\left|\psi_{\perp}\right\rangle$ is unknown, and presumed not orthogonal to any of the energy eigenkets. We can still calculate the norm of the trial function

$$
\begin{align*}
\langle\phi \mid \phi\rangle & =\sum_{n, m}\left\langle c_{n} \psi_{n}+c_{\perp} \psi_{\perp} \mid c_{m} \psi_{m}+c_{\perp} \psi_{\perp}\right\rangle \\
& =\sum_{n}\left|c_{n}\right|^{2}+c_{n}^{*} c_{\perp}\left\langle\psi_{n} \mid \psi_{\perp}\right\rangle+c_{n} c_{\perp}^{*}\left\langle\psi_{\perp} \mid \psi_{n}\right\rangle+\left|c_{\perp}\right|^{2}\left\langle\psi_{\perp} \mid \psi_{\perp}\right\rangle  \tag{2.84}\\
& =\left\langle\psi_{\perp} \mid \psi_{\perp}\right\rangle+\sum_{n}\left|c_{n}\right|^{2}+2 \operatorname{Re}\left(c_{n}^{*} c_{\perp}\left\langle\psi_{n} \mid \psi_{\perp}\right\rangle\right) .
\end{align*}
$$

Similarly we can calculate the energy expectation for this unnormalized state and find

$$
\begin{align*}
\langle\phi| H|\phi\rangle= & \sum_{n, m}\left\langle c_{n} \psi_{n}+c_{\perp} \psi_{\perp}\right| H\left|c_{m} \psi_{m}+c_{\perp} \psi_{\perp}\right\rangle \\
= & \sum_{n}\left|c_{n}\right|^{2} E_{n}+c_{n}^{*} c_{\perp} E_{n}\left\langle\psi_{n} \mid \psi_{\perp}\right\rangle  \tag{2.85}\\
& +c_{n} c_{\perp}^{*} E_{n}\left\langle\psi_{\perp} \mid \psi_{n}\right\rangle+\left|c_{\perp}\right|^{2}\left\langle\psi_{\perp}\right| H\left|\psi_{\perp}\right\rangle
\end{align*}
$$

Our normalized energy expectation is therefore the considerably messier

$$
\begin{align*}
& \frac{\langle\phi| H|\phi\rangle}{\langle\phi \mid \phi\rangle} \\
& =\frac{\sum_{n}\left|c_{n}\right|^{2} E_{n}+c_{n}^{*} c_{\perp} E_{n}\left\langle\psi_{n} \mid \psi_{\perp}\right\rangle+c_{n} c_{\perp}^{*} E_{n}\left\langle\psi_{\perp} \mid \psi_{n}\right\rangle+\left|c_{\perp}\right|^{2}\left\langle\psi_{\perp}\right| H\left|\psi_{\perp}\right\rangle}{\left\langle\psi_{\perp} \mid \psi_{\perp}\right\rangle+\sum_{m}\left|c_{m}\right|^{2}+2 \operatorname{Re}\left(c_{m}^{*} c_{\perp}\left\langle\psi_{m} \mid \psi_{\perp}\right\rangle\right)} \\
& \geq \frac{\sum_{n}\left|c_{n}\right|^{2} E_{0}+c_{n}^{*} c_{\perp} E_{n}\left\langle\psi_{n} \mid \psi_{\perp}\right\rangle+c_{n} c_{\perp}^{*} E_{n}\left\langle\psi_{\perp} \mid \psi_{n}\right\rangle+\left|c_{\perp}\right|^{2}\left\langle\psi_{\perp}\right| H\left|\psi_{\perp}\right\rangle}{\left\langle\psi_{\perp} \mid \psi_{\perp}\right\rangle+\sum_{m}\left|c_{m}\right|^{2}+2 \operatorname{Re}\left(c_{m}^{*} c_{\perp}\left\langle\psi_{m} \mid \psi_{\perp}\right\rangle\right)} . \tag{2.86}
\end{align*}
$$

With a requirement to include the perpendicular cross terms the norm does not just cancel out, leaving us with a clean estimation of the ground state energy. In order to utilize this variational method, we implicitly have an assumption that the $\left\langle\psi_{\perp} \mid \psi_{\perp}\right\rangle$ and $\left\langle\psi_{m} \mid \psi_{\perp}\right\rangle$ terms in the denominator are sufficiently small that they can be neglected.

Calculating the Fourier terms. In order to see how much a problem representing this trial function in the Harmonic oscillator wavefunction solution space, we can just calculate the Fourier fit.

Our first few basis functions, with $\alpha=\sqrt{m \omega / \hbar}$ are

$$
\begin{align*}
& u_{0}=\sqrt{\frac{\alpha}{\sqrt{\pi}}} e^{-\alpha^{2} x^{2} / 2} \\
& u_{1}=\sqrt{\frac{\alpha}{2 \sqrt{\pi}}}(2 \alpha x) e^{-\alpha^{2} x^{2} / 2}  \tag{2.87}\\
& u_{2}=\sqrt{\frac{\alpha}{8 \sqrt{\pi}}}\left(4 \alpha^{2} x^{2}-2\right) e^{-\alpha^{2} x^{2} / 2} .
\end{align*}
$$

In general our wavefunctions are

$$
\begin{align*}
u_{n} & =N_{n} H_{n}(\alpha x) e^{-\alpha^{2} x^{2} / 2} \\
N_{n} & =\sqrt{\frac{\alpha}{\sqrt{\pi} 2^{n} n!}}  \tag{2.88}\\
H_{n}(\eta) & =(-1)^{n} e^{\eta^{2}} \frac{d^{n}}{d \eta^{n}} e^{-\eta^{2}},
\end{align*}
$$

from which we find

$$
\begin{equation*}
\psi(x)=e^{-\alpha^{2} x^{2} / 2}\left(N_{n}\right)^{2} H_{n}(\alpha x) \int_{-\infty}^{\infty} H_{n}(\alpha x) e^{-\alpha^{2} x^{2} / 2} \psi(x) d x \tag{2.89}
\end{equation*}
$$

Our wave function, with $\beta=1$ is plotted in fig. 2.3. The zeroth order


Figure 2.3: Exponential trial function with absolute exponential die off.
fitting using the Gaussian exponential is found to be

$$
\begin{equation*}
\psi_{0}(x)=\sqrt{2 \beta} \operatorname{erfc}\left(\frac{\beta}{\sqrt{2} \alpha}\right) e^{-\alpha^{2} x^{2} / 2+\beta^{2} /\left(2 \alpha^{2}\right)} \tag{2.90}
\end{equation*}
$$

With $\alpha=\beta=1$, this is plotted in fig. 2.4 and can be seen to match fairly well The higher order terms get small fast, but we can see in fig. 2.5, where a tenth order fitting is depicted that it would take a number of them


Figure 2.4: First ten orders, fitting harmonic oscillator wavefunctions to this trial function.


Figure 2.5: Tenth order harmonic oscillator wavefunction fitting.
to get anything close to the sharp peak that we have in our exponential trial function. Note that all the brakets of even orders in $n$ with the trial function are zero, which is why the tenth order approximation is only a sum of six terms.

Details for this harmonic oscillator wavefunction fitting can be found in gaussian_fitting_for_abs_function.nb can be found separately, calculated using a Mathematica worksheet.

The question of interest is why we can approximate the trial function so nicely (except at the origin) even with just a first order approximation (polynomial times Gaussian functions where the polynomials are Hankel functions), and we can get an exact value for the lowest energy state using the first order approximation of our trial function, why do we get garbage from the variational method, where enough terms are implicitly included that the peak should be sharp. It must therefore be important to consider the origin, but how do we give some meaning to the derivative of the absolute value function? The key (supplied when asking Professor Sipe in office hours for the course) is to express the absolute value function in terms of Heavyside step functions, for which the derivative can be identified as the delta function.

Correcting, treating the origin this way. Here is how we can express the absolute value function using the Heavyside step

$$
\begin{equation*}
|x|=x \theta(x)-x \theta(-x), \tag{2.91}
\end{equation*}
$$

where the step function is zero for $x<0$ and one for $x>0$ as plotted in fig. 2.6. Expressed this way, with the identification $\theta^{\prime}(x)=\delta(x)$, we have


Figure 2.6: Unit step function.
for the derivative of the absolute value function

$$
\begin{align*}
|x|^{\prime} & =x^{\prime} \theta(x)-x^{\prime} \theta(-x)+x \theta^{\prime}(x)-x \theta^{\prime}(-x) \\
& =\theta(x)-\theta(-x)+x \delta(x)+x \delta(-x)  \tag{2.92}\\
& =\theta(x)-\theta(-x)+x \delta(x)+x \delta(x) .
\end{align*}
$$

Observe that we have our expected unit derivative for $x>0$, and -1 derivative for $x<0$. At the origin our $\theta$ contributions vanish, and we are left with

$$
\begin{equation*}
\left.|x|^{\prime}\right|_{x=0}=\left.2 x \delta(x)\right|_{x=0} . \tag{2.93}
\end{equation*}
$$

We have got zero times infinity here, so how do we give meaning to this? As with any delta functional, we have got to apply it to a well behaved (square integrable) test function $f(x)$ and integrate. Doing so we have

$$
\begin{align*}
\int_{-\infty}^{\infty} d x|x|^{\prime} f(x) & =2 \int_{-\infty}^{\infty} d x x \delta(x) f(x)  \tag{2.94}\\
& =2(0) f(0) .
\end{align*}
$$

This equals zero for any well behaved test function $f(x)$. Since the delta function only picks up the contribution at the origin, we can therefore identify $|x|^{\prime}$ as zero at the origin.

Using the same technique, we can express our trial function in terms of steps

$$
\begin{equation*}
\psi=e^{-\beta|x|}=\theta(x) e^{-\beta x}+\theta(-x) e^{\beta x} . \tag{2.95}
\end{equation*}
$$

This we can now take derivatives of, even at the origin, and find

$$
\begin{align*}
\psi^{\prime} & =\theta^{\prime}(x) e^{-\beta x}+\theta^{\prime}(-x) e^{\beta x}-\beta \theta(x) e^{-\beta x}+\beta \theta(-x) e^{\beta x} \\
& =\delta(x) e^{-\beta x}-\delta(-x) e^{\beta x}-\beta \theta(x) e^{-\beta x}+\beta \theta(-x) e^{\beta x} \\
& =\delta(x) e^{-\beta x}-\delta(x) e^{\beta x}-\beta \theta(x) e^{-\beta x}+\beta \theta(-x) e^{\beta x}  \tag{2.96}\\
& =\beta\left(-\theta(x) e^{-\beta x}+\theta(-x) e^{\beta x}\right) .
\end{align*}
$$

Taking second derivatives we find

$$
\begin{align*}
\psi^{\prime \prime} & =\beta\left(-\theta^{\prime}(x) e^{-\beta x}+\theta^{\prime}(-x) e^{\beta x}+\beta \theta(x) e^{-\beta x}+\beta \theta(-x) e^{\beta x}\right) \\
& =\beta\left(-\delta(x) e^{-\beta x}-\delta(-x) e^{\beta x}+\beta \theta(x) e^{-\beta x}+\beta \theta(-x) e^{\beta x}\right)  \tag{2.97}\\
& =\beta^{2} \psi-2 \beta \delta(x) .
\end{align*}
$$

Now application of the Hamiltonian operator on our trial function gives us

$$
\begin{equation*}
H \psi=-\frac{\hbar^{2}}{2 m}\left(\beta^{2} \psi-2 \delta(x)\right)+\frac{1}{2} m \omega^{2} x^{2} \psi \tag{2.98}
\end{equation*}
$$

So

$$
\begin{align*}
\langle\psi| H|\psi\rangle & =\int_{-\infty}^{\infty}\left(-\frac{\hbar^{2} \beta^{2}}{2 m}+\frac{1}{2} m \omega^{2} x^{2}\right) e^{-2 \beta|x|}+\frac{\hbar^{2} \beta}{m} \int_{-\infty}^{\infty} \delta(x) e^{-\beta|x|} \\
& =-\frac{\beta \hbar^{2}}{2 m}+\frac{m \omega^{2}}{4 \beta^{3}}+\frac{\hbar^{2} \beta}{m}  \tag{2.99}\\
& =\frac{\beta \hbar^{2}}{2 m}+\frac{m \omega^{2}}{4 \beta^{3}}
\end{align*}
$$

Normalized we have

$$
\begin{equation*}
E[\beta]=\frac{\langle\psi| H|\psi\rangle}{\langle\psi \mid \psi\rangle}=\frac{\beta^{2} \hbar^{2}}{2 m}+\frac{m \omega^{2}}{4 \beta^{2}} . \tag{2.100}
\end{equation*}
$$

This is looking much more promising. We will have the sign alternation that we require to find a positive, non-complex, value for $\beta$ when $E[\beta]$ is minimized. That is

$$
\begin{equation*}
0=\frac{\partial E}{\partial \beta}=\frac{\beta \hbar^{2}}{m}-\frac{m \omega^{2}}{2 \beta^{3}} \tag{2.101}
\end{equation*}
$$

so the extremum is found at

$$
\begin{equation*}
\beta^{4}=\frac{m^{2} \omega^{2}}{2 \hbar^{2}} \tag{2.102}
\end{equation*}
$$

Plugging this back in we find that our trial function associated with the minimum energy (unnormalized still) is

$$
\begin{equation*}
\psi=e^{-\sqrt{\frac{m \omega x^{2}}{\sqrt{2} \hbar}}} \tag{2.103}
\end{equation*}
$$

and that energy, after substitution, is

$$
\begin{equation*}
E\left[\beta_{\min }\right]=\frac{\hbar \omega}{2} \sqrt{2} \tag{2.104}
\end{equation*}
$$

We have something that is $1.4 \times$ the true ground state energy, but is at least a ball park value. However, to get this result, we have to be very careful to treat our point of singularity. A derivative that we would call undefined in first year calculus, is not only defined, but required, for this treatment to work!

See $\S 16.1$ of the text [4].
We can sometimes use this sort of physical insight to help construct a good approximation. This is provided that we have some of this physical insight, or that it is good insight in the first place. This is the no-think (turn the crank) approach. Here we split our Hamiltonian into two parts

$$
\begin{equation*}
H=H_{0}+H^{\prime} \tag{3.1}
\end{equation*}
$$

where $H_{0}$ is a Hamiltonian for which we know the energy eigenstates and the eigenkets. The $H^{\prime}$ is the "perturbation" that is supposed to be small "in some sense".

Prof Sine will provide some references later that provide a more seecific meaning to this "smallness". From some ad-hoc discussion in the class it sounds like one has to consider sequences of operators, and look at the convergence of those sequences (is this L2 measure theory?) We


Figure 3.1: Example of small perturbation from known Hamiltonian.
would like to consider a range of problems of the form

$$
\begin{equation*}
H=H_{0}+\lambda H^{\prime}, \tag{3.2}
\end{equation*}
$$

where

$$
\begin{equation*}
\lambda \in[0,1] . \tag{3.3}
\end{equation*}
$$

When $\lambda \rightarrow 0$ we have

$$
\begin{equation*}
H \rightarrow H_{0} \tag{3.4}
\end{equation*}
$$

which is the problem that we already know. However, for $\lambda \rightarrow 1$ we have

$$
\begin{equation*}
H=H_{0}+H^{\prime} \tag{3.5}
\end{equation*}
$$

the problem that we would like to solve. We are assuming that we know the eigenstates and eigenvalues for $H_{0}$. Assuming no degeneracy

$$
\begin{equation*}
H_{0}\left|\psi_{s}^{(0)}\right\rangle=E_{s}^{(0)}\left|\psi_{s}^{(0)}\right\rangle \tag{3.6}
\end{equation*}
$$

We seek

$$
\begin{equation*}
\left(H_{0}+H^{\prime}\right)\left|\psi_{s}\right\rangle=E_{s}\left|\psi_{s}\right\rangle \tag{3.7}
\end{equation*}
$$

which is the $\lambda=1$ case. Once (if) found, when $\lambda \rightarrow 0$ we will have

$$
\begin{align*}
E_{s} & \rightarrow E_{s}^{(0)} \\
\left|\psi_{s}\right\rangle & \rightarrow\left|\psi_{s}^{(0)}\right\rangle \tag{3.8}
\end{align*}
$$

We expand the energy as

$$
\begin{equation*}
E_{s}=E_{s}^{(0)}+\lambda E_{s}^{(1)}+\lambda^{2} E_{s}^{(2)} \tag{3.9}
\end{equation*}
$$

where

$$
\begin{equation*}
\psi_{s}=\sum_{n} c_{n s}\left|\psi_{n}^{(0)}\right\rangle \tag{3.10}
\end{equation*}
$$

This we know we can do because we are assumed to have a complete set of states. With

$$
\begin{equation*}
c_{n s}=c_{n s}^{(0)}+\lambda c_{n s}^{(1)}+\lambda^{2} c_{n s}^{(2)} \tag{3.11}
\end{equation*}
$$

where

$$
\begin{equation*}
c_{n s}^{(0)}=\delta_{n s} \tag{3.12}
\end{equation*}
$$

There is a subtlety here that will be treated differently from the text. We write

$$
\begin{align*}
\left|\psi_{s}\right\rangle & =\left|\psi_{s}^{(0)}\right\rangle+\lambda \sum_{n} c_{n s}^{(1)}\left|\psi_{n}^{(0)}\right\rangle+\lambda^{2} \sum_{n} c_{n s}^{(2)}\left|\psi_{n}^{(0)}\right\rangle+\cdots \\
& =\left(1+\lambda c_{s s}^{(1)}+\cdots\right)\left|\psi_{s}^{(0)}\right\rangle+\lambda \sum_{n \neq s} c_{n s}^{(1)}\left|\psi_{n}^{(0)}\right\rangle+\cdots \tag{3.13}
\end{align*}
$$

Take

$$
\begin{align*}
\left|\bar{\psi}_{s}\right\rangle & =\left|\bar{\psi}_{s}^{(0)}\right\rangle+\lambda \frac{\sum_{n \neq s} c_{n s}^{(1)}\left|\psi_{n}^{(0)}\right\rangle}{1+\lambda c_{s s}^{(1)}}+\cdots  \tag{3.14}\\
& =\left|\bar{\psi}_{s}^{(0)}\right\rangle+\lambda \sum_{n \neq s} \bar{c}_{n s}^{(1)}\left|\psi_{n}^{(0)}\right\rangle+\cdots
\end{align*}
$$

where

$$
\begin{equation*}
\bar{c}_{n s}^{(1)}=\frac{c_{n s}^{(1)}}{1+\lambda c_{s s}^{(1)}} \tag{3.15}
\end{equation*}
$$

We have:

$$
\begin{align*}
& \bar{c}_{n s}^{(1)}=c_{n s}^{(1)} \\
& \bar{c}_{n s}^{(2)} \neq c_{n s}^{(2)} . \tag{3.16}
\end{align*}
$$

FIXME: I missed something here.
Note that this is no longer normalized

$$
\begin{equation*}
\left\langle\bar{\psi}_{s} \mid \bar{\psi}_{s}\right\rangle \neq 1 \tag{3.17}
\end{equation*}
$$

### 3.1 TIME INDEPENDENT PERTURBATION.

The setup. To recap, we were covering the time independent perturbation methods from §16.1 of the text [4]. We start with a known Hamiltonian $H_{0}$, and alter it with the addition of a "small" perturbation

$$
\begin{equation*}
H=H_{0}+\lambda H^{\prime}, \quad \lambda \in[0,1] \tag{3.18}
\end{equation*}
$$

For the original operator, we assume that a complete set of eigenvectors and eigenkets is known

$$
\begin{equation*}
H_{0}\left|\psi_{s}{ }^{(0)}\right\rangle=E_{s}^{(0)}\left|\psi_{s}^{(0)}\right\rangle \tag{3.19}
\end{equation*}
$$

We seek the perturbed eigensolution

$$
\begin{equation*}
H\left|\psi_{s}\right\rangle=E_{s}\left|\psi_{s}\right\rangle, \tag{3.20}
\end{equation*}
$$

and assumed a perturbative series representation for the energy eigenvalues in the new system

$$
\begin{equation*}
E_{s}=E_{s}^{(0)}+\lambda E_{s}^{(1)}+\lambda^{2} E_{s}^{(2)}+\cdots \tag{3.21}
\end{equation*}
$$

Given an assumed representation for the new eigenkets in terms of the known basis

$$
\begin{equation*}
\left|\psi_{s}\right\rangle=\sum_{n} c_{n s}\left|\psi_{n}^{(0)}\right\rangle \tag{3.22}
\end{equation*}
$$

and a perturbative series representation for the probability coefficients

$$
\begin{equation*}
c_{n s}=c_{n s}^{(0)}+\lambda c_{n s}^{(1)}+\lambda^{2} c_{n s}^{(2)} \tag{3.23}
\end{equation*}
$$

so that

$$
\begin{equation*}
\left|\psi_{s}\right\rangle=\sum_{n} c_{n s}^{(0)}\left|\psi_{n}^{(0)}\right\rangle+\lambda \sum_{n} c_{n s}^{(1)}\left|\psi_{n}^{(0)}\right\rangle+\lambda^{2} \sum_{n} c_{n s}^{(2)}\left|\psi_{n}^{(0)}\right\rangle+\cdots \tag{3.24}
\end{equation*}
$$

Setting $\lambda=0$ requires

$$
\begin{equation*}
c_{n s}{ }^{(0)}=\delta_{n s} \tag{3.25}
\end{equation*}
$$

for

$$
\begin{align*}
\left|\psi_{s}\right\rangle= & \left|\psi_{s}^{(0)}\right\rangle+\lambda \sum_{n} c_{n s}{ }^{(1)}\left|\psi_{n}^{(0)}\right\rangle+\lambda^{2} \sum_{n} c_{n s}^{(2)}\left|\psi_{n}^{(0)}\right\rangle+\cdots \\
= & \left(1+\lambda c_{s s}{ }^{(1)}+\lambda^{2} c_{s s}{ }^{(2)}+\cdots\right)\left|\psi_{s}^{(0)}\right\rangle  \tag{3.26}\\
& +\lambda \sum_{n \neq s} c_{n s}{ }^{(1)}\left|\psi_{n}^{(0)}\right\rangle+\lambda^{2} \sum_{n \neq s} c_{n s}^{(2)}\left|\psi_{n}^{(0)}\right\rangle+\cdots
\end{align*}
$$

We rescaled our kets

$$
\begin{equation*}
\left|\bar{\psi}_{s}\right\rangle=\left|\psi_{s}^{(0)}\right\rangle+\lambda \sum_{n \neq s} \bar{c}_{n s}^{(1)}\left|\psi_{n}^{(0)}\right\rangle+\lambda^{2} \sum_{n \neq s} \bar{c}_{n s}^{(2)}\left|\psi_{n}^{(0)}\right\rangle+\cdots \tag{3.27}
\end{equation*}
$$

where

The normalization of the rescaled kets is then

$$
\begin{equation*}
\left\langle\bar{\psi}_{s} \mid \bar{\psi}_{s}\right\rangle=1+\lambda^{2} \sum_{n \neq s}\left|\bar{c}_{s s}^{(1)}\right|^{2}+\cdots \equiv \frac{1}{Z_{s}} \tag{3.29}
\end{equation*}
$$

One can then construct a renormalized ket if desired

$$
\begin{equation*}
\left|\bar{\psi}_{s}\right\rangle_{R}=Z_{s}^{1 / 2}\left|\bar{\psi}_{s}\right\rangle \tag{3.30}
\end{equation*}
$$

so that

$$
\begin{equation*}
\left(\left|\bar{\psi}_{s}\right\rangle_{R}\right)^{\dagger}\left|\bar{\psi}_{s}\right\rangle_{R}=Z_{s}\left\langle\bar{\psi}_{s} \mid \bar{\psi}_{s}\right\rangle=1 \tag{3.31}
\end{equation*}
$$

The meat. $\quad$ That is as far as we got last time. We continue by renaming terms in eq. (3.27)

$$
\begin{equation*}
\left|\bar{\psi}_{s}\right\rangle=\left|\psi_{s}^{(0)}\right\rangle+\lambda\left|\psi_{s}^{(1)}\right\rangle+\lambda^{2}\left|\psi_{s}^{(2)}\right\rangle+\cdots, \tag{3.32}
\end{equation*}
$$

where

$$
\begin{equation*}
\left|\psi_{s}{ }^{(j)}\right\rangle=\sum_{n \neq s} \bar{c}_{n s}{ }^{(j)}\left|\psi_{n}{ }^{(0)}\right\rangle \tag{3.33}
\end{equation*}
$$

Now we act on this with the Hamiltonian

$$
\begin{equation*}
H\left|\bar{\psi}_{s}\right\rangle=E_{s}\left|\bar{\psi}_{s}\right\rangle \tag{3.34}
\end{equation*}
$$

or

$$
\begin{equation*}
H\left|\bar{\psi}_{s}\right\rangle-E_{s}\left|\bar{\psi}_{s}\right\rangle=0 \tag{3.35}
\end{equation*}
$$

Expanding this, we have

$$
\begin{align*}
0= & \left(H_{0}+\lambda H^{\prime}\right)\left(\left|\psi_{s}^{(0)}\right\rangle+\lambda\left|\psi_{s}^{(1)}\right\rangle+\lambda^{2}\left|\psi_{s}^{(2)}\right\rangle+\cdots\right) \\
& -\left(E_{s}^{(0)}+\lambda E_{s}^{(1)}+\lambda^{2} E_{s}^{(2)}+\cdots\right) \times  \tag{3.36}\\
& \left(\left|\psi_{s}^{(0)}\right\rangle+\lambda\left|\psi_{s}^{(1)}\right\rangle+\lambda^{2}\left|\psi_{s}^{(2)}\right\rangle+\cdots\right)=0 .
\end{align*}
$$

We want to write this as

$$
\begin{equation*}
|A\rangle+\lambda|B\rangle+\lambda^{2}|C\rangle+\cdots=0 \tag{3.37}
\end{equation*}
$$

This is

$$
\begin{align*}
0 & =\lambda^{0}\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s}^{(0)}\right\rangle \\
& +\lambda\left(\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s}^{(1)}\right\rangle+\left(H^{\prime}-E_{s}^{(1)}\right)\left|\psi_{s}^{(0)}\right\rangle\right) \\
& +\lambda^{2}\left(\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s}^{(2)}\right\rangle+\left(H^{\prime}-E_{s}^{(1)}\right)\left|\psi_{s}^{(1)}\right\rangle-E_{s}^{(2)}\left|\psi_{s}^{(0)}\right\rangle\right) \tag{3.38}
\end{align*}
$$

So we form

$$
\begin{align*}
|A\rangle & =\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s}^{(0)}\right\rangle \\
|B\rangle & =\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s}^{(1)}\right\rangle+\left(H^{\prime}-E_{s}^{(1)}\right)\left|\psi_{s}^{(0)}\right\rangle  \tag{3.39}\\
|C\rangle & =\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s}^{(2)}\right\rangle+\left(H^{\prime}-E_{s}^{(1)}\right)\left|\psi_{s}^{(1)}\right\rangle-E_{s}^{(2)}\left|\psi_{s}^{(0)}\right\rangle
\end{align*}
$$

and so forth.

Zeroth order in $\lambda . \quad$ Since $H_{0}\left|\psi_{s}{ }^{(0)}\right\rangle=E_{s}^{(0)}\left|\psi_{s}{ }^{(0)}\right\rangle$, this first condition on $|A\rangle$ is not much more than a statement that $0-0=0$.

First order in $\lambda$. How about $|B\rangle=0$ ? For this to be zero we require that both of the following are simultaneously zero

$$
\begin{align*}
\left\langle\psi_{s}^{(0)} \mid B\right\rangle & =0 \\
\left\langle\psi_{m}^{(0)} \mid B\right\rangle & =0, \quad m \neq s \tag{3.40}
\end{align*}
$$

This first condition is

$$
\begin{equation*}
\left\langle\psi_{s}^{(0)}\right|\left(H^{\prime}-E_{s}^{(1)}\right)\left|\psi_{s}^{(0)}\right\rangle=0 \tag{3.41}
\end{equation*}
$$

With

$$
\begin{equation*}
\left\langle\psi_{m}^{(0)}\right| H^{\prime}\left|\psi_{s}^{(0)}\right\rangle \equiv H_{m s}^{\prime} \tag{3.42}
\end{equation*}
$$

or

$$
\begin{equation*}
H_{s s}^{\prime}=E_{s}^{(1)} \tag{3.43}
\end{equation*}
$$

From the second condition we have

$$
\begin{equation*}
0=\left\langle\psi_{m}^{(0)}\right|\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s}^{(1)}\right\rangle+\left\langle\psi_{m}^{(0)}\right|\left(H^{\prime}-E_{s}^{(1)}\right)\left|\psi_{s}^{(0)}\right\rangle \tag{3.44}
\end{equation*}
$$

Utilizing the Hermitian nature of $H_{0}$ we can act backwards on $\left\langle\psi_{m}{ }^{(0)}\right|$

$$
\begin{equation*}
\left\langle\psi_{m}^{(0)}\right| H_{0}=E_{m}^{(0)}\left\langle\psi_{m}^{(0)}\right| \tag{3.45}
\end{equation*}
$$

We note that $\left\langle\psi_{m}{ }^{(0)} \mid \psi_{s}{ }^{(0)}\right\rangle=0, m \neq s$. We can also expand the $\left\langle\psi_{m}{ }^{(0)} \mid \psi_{s}{ }^{(1)}\right\rangle$, which is

$$
\begin{equation*}
\left\langle\psi_{m}{ }^{(0)} \mid \psi_{s}^{(1)}\right\rangle=\left\langle\psi_{m}^{(0)}\right|\left(\sum_{n \neq s} \bar{c}_{n s}^{(1)}\left|\psi_{n}^{(0)}\right\rangle\right) \tag{3.46}
\end{equation*}
$$

I found that reducing this sum was not obvious until some actual integers were plugged in. Suppose that $s=3$, and $m=5$, then this is

$$
\begin{align*}
\left\langle\psi_{5}{ }^{(0)} \mid \psi_{3}{ }^{(1)}\right\rangle & \left.=\left\langle\psi_{5}{ }^{(0)}\right| \sum_{n=0,1,2,4,5, \cdots} \bar{c}_{n 3}{ }^{(1)}\left|\psi_{n}{ }^{(0)}\right\rangle\right) \\
& =\bar{c}_{53}{ }^{(1)}\left\langle\psi_{5}{ }^{(0)} \mid \psi_{5}{ }^{(0)}\right\rangle  \tag{3.47}\\
& =\bar{c}_{53}{ }^{(1)} .
\end{align*}
$$

Observe that we can also replace the superscript (1) with $(j)$ in the above manipulation without impacting anything else. That and putting back in the abstract indices, we have the general result

$$
\begin{equation*}
\left\langle\psi_{m}{ }^{(0)} \mid \psi_{s}{ }^{(j)}\right\rangle=\bar{c}_{m s}{ }^{(j)} \tag{3.48}
\end{equation*}
$$

Utilizing this gives us

$$
\begin{equation*}
0=\left(E_{m}^{(0)}-E_{s}^{(0)}\right) \bar{c}_{m s}^{(1)}+H_{m s}^{\prime} . \tag{3.49}
\end{equation*}
$$

Summarizing what we learn from our $|B\rangle=0$ conditions we have

$$
\begin{align*}
E_{s}^{(1)} & ={H^{\prime}}_{s s} \\
\bar{c}_{m s}{ }^{(1)} & =\frac{{H^{\prime}{ }_{m s}}_{E_{s}^{(0)}-E_{m}^{(0)}}}{} . \tag{3.50}
\end{align*}
$$

Second order in $\lambda$. Doing the same thing for $|C\rangle=0$ we form (or assume)

$$
\begin{equation*}
\left\langle\psi_{s}^{(0)} \mid C\right\rangle=0 \tag{3.51}
\end{equation*}
$$

and

$$
\begin{align*}
0= & \left\langle\psi_{s}^{(0)} \mid C\right\rangle \\
= & \left\langle\psi_{s}^{(0)}\right|\left(\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s}^{(2)}\right\rangle+\left(H^{\prime}-E_{s}^{(1)}\right)\left|\psi_{s}^{(1)}\right\rangle-E_{s}^{(2)}\left|\psi_{s}^{(0)}\right\rangle\right)  \tag{3.52}\\
= & \left(E_{s}^{(0)}-E_{s}^{(0)}\right)\left\langle\psi_{s}^{(0)} \mid \psi_{s}^{(2)}\right\rangle \\
& +\left\langle\psi_{s}^{(0)}\right|\left(H^{\prime}-E_{s}^{(1)}\right)\left|\psi_{s}^{(1)}\right\rangle-E_{s}^{(2)}\left\langle\psi_{s}^{(0)} \mid \psi_{s}^{(0)}\right\rangle .
\end{align*}
$$

We need to know what the $\left\langle\psi_{s}{ }^{(0)} \mid \psi_{s}{ }^{(1)}\right\rangle$ is, and find that it is zero

$$
\begin{equation*}
\left\langle\psi_{s}^{(0)} \mid \psi_{s}^{(1)}\right\rangle=\left\langle\psi_{s}^{(0)}\right| \sum_{n \neq s} \bar{c}_{n s}^{(1)}\left|\psi_{n}^{(0)}\right\rangle . \tag{3.53}
\end{equation*}
$$

Again, suppose that $s=3$. Our sum ranges over all $n \neq 3$, so all the brakets are zero. Utilizing that we have

$$
\begin{align*}
E_{s}^{(2)} & =\left\langle\psi_{s}^{(0)}\right| H^{\prime}\left|\psi_{s}^{(1)}\right\rangle \\
& =\left\langle\psi_{s}^{(0)}\right| H^{\prime} \sum_{m \neq s} \bar{c}_{m s}{ }^{(1)}\left|\psi_{m}^{(0)}\right\rangle  \tag{3.54}\\
& =\sum_{m \neq s} \bar{c}_{m s}{ }^{(1)} H_{s m}^{\prime} .
\end{align*}
$$

From eq. (3.50) we have

$$
\begin{equation*}
E_{s}^{(2)}=\sum_{m \neq s} \frac{H^{\prime}{ }_{m s}}{E_{s}^{(0)}-E_{m}^{(0)}} H^{\prime}{ }_{s m}=\sum_{m \neq s} \frac{\left|H^{\prime}{ }_{m s}\right|^{2}}{E_{s}^{(0)}-E_{m}^{(0)}} \tag{3.55}
\end{equation*}
$$

We can now summarize by forming the first order terms of the perturbed energy and the corresponding kets

$$
\begin{align*}
E_{s} & =E_{s}^{(0)}+\lambda H_{s s}^{\prime}+\lambda^{2} \sum_{m \neq s} \frac{\left|H_{m s}^{\prime}\right|^{2}}{E_{s}^{(0)}-E_{m}^{(0)}}+\cdots  \tag{3.56}\\
\left|\bar{\psi}_{s}\right\rangle & =\left|\psi_{s}^{(0)}\right\rangle+\lambda \sum_{m \neq s} \frac{H_{m s}^{\prime}}{E_{s}^{(0)}-E_{m}^{(0)}}\left|\psi_{m}^{(0)}\right\rangle+\cdots
\end{align*}
$$

We can continue calculating, but are hopeful that we can stop the calculation without doing more work, even if $\lambda=1$. If one supposes that the

$$
\begin{equation*}
\sum_{m \neq s} \frac{H^{\prime}{ }_{m s}}{E_{s}^{(0)}-E_{m}^{(0)}}, \tag{3.57}
\end{equation*}
$$

term is "small", then we can hope that truncating the sum will be reasonable for $\lambda=1$. This would be the case if

$$
\begin{equation*}
H_{m s}^{\prime} \ll\left|E_{s}^{(0)}-E_{m}^{(0)}\right|, \tag{3.58}
\end{equation*}
$$

however, to put some mathematical rigour into making a statement of such smallness takes a lot of work. We are referred to [12]. Incidentally, these are loosely referred to as the first and second testaments, because of the author's name, and the fact that they came as two volumes historically.

### 3.2 ISSUES CONCERNING DEGENERACY.

When the perturbed state is non-degenerate. Suppose the state of interest is non-degenerate but others are

FIXME: diagram. states designated by dashes labeled $n 1, n 2, n 3$ degeneracy $\alpha=3$ for energy $E_{n}^{(0)}$.

This is no problem except for notation, and if the analysis is repeated we find

$$
\begin{gather*}
E_{s}=E_{s}^{(0)}+\lambda H^{\prime}{ }_{s s}+\lambda^{2} \sum_{m \neq s, \alpha} \frac{\left|H_{m \alpha ; s}^{\prime}\right|^{2}}{E_{s}^{(0)}-E_{m}^{(0)}}+\cdots  \tag{3.59}\\
\left|\bar{\psi}_{s}\right\rangle=\left|\psi_{s}{ }^{(0)}\right\rangle+\lambda \sum_{m \neq s, \alpha} \frac{H^{\prime}{ }_{m \alpha ; s}}{E_{s}^{(0)}-E_{m}^{(0)}}\left|\psi_{m \alpha}{ }^{(0)}\right\rangle+\cdots, \tag{3.60}
\end{gather*}
$$

where

$$
\begin{equation*}
H_{m \alpha ; s}^{\prime}=\left\langle\psi_{m \alpha}{ }^{(0)}\right| H^{\prime}\left|\psi_{s \alpha}{ }^{(0)}\right\rangle . \tag{3.61}
\end{equation*}
$$

When the perturbed state is also degenerate. FIXME: diagram. states designated by dashes labeled $n 1, n 2, n 3$ degeneracy $\alpha=3$ for energy $E_{n}^{(0)}$, and states designated by dashes labeled $s 1, s 2, s 3$ degeneracy $\alpha=3$ for energy $E_{s}^{(0)}$.

If we just blindly repeat the derivation for the non-degenerate case we would obtain

$$
\begin{align*}
E_{s}= & E_{s}^{(0)}+\lambda H_{s 1 ; s 1}^{\prime}+\lambda^{2} \sum_{m \neq s, \alpha} \frac{\left|H^{\prime}{ }_{m \alpha ; s 1}\right|^{2}}{E_{s}^{(0)}-E_{m}^{(0)}}  \tag{3.62}\\
& +\lambda^{2} \sum_{\alpha \neq 1} \frac{\left|H_{s \alpha ; s 1}^{\prime}\right|^{2}}{E_{s}^{(0)}-E_{s}^{(0)}}+\cdots  \tag{3.63}\\
\left|\bar{\psi}_{s}\right\rangle= & \left|\psi_{s}^{(0)}\right\rangle+\lambda \sum_{m \neq s, \alpha} \frac{H_{m \alpha ; s}^{\prime}}{E_{s}^{(0)}-E_{m}^{(0)}}\left|\psi_{m \alpha}{ }^{(0)}\right\rangle  \tag{3.64}\\
& +\lambda \sum_{\alpha \neq s 1} \frac{H_{s \alpha ; s 1}^{\prime(0)}\left|\psi_{s \alpha}^{(0)}\right\rangle+\cdots,}{E_{s}^{(0)}}{ }^{(0)}, \tag{3.65}
\end{align*}
$$

where

$$
\begin{equation*}
H_{m \alpha ; s 1}^{\prime}=\left\langle\psi_{m \alpha}{ }^{(0)}\right| H^{\prime}\left|\psi_{s 1}{ }^{(0)}\right\rangle . \tag{3.66}
\end{equation*}
$$

Note that the $E_{s}^{(0)}-E_{s}^{(0)}$ is NOT a typo, and why we run into trouble. There is one case where a perturbation approach is still possible. That case is if we happen to have

$$
\begin{equation*}
\left\langle\psi_{m a}{ }^{(0)}\right| H^{\prime}\left|\psi_{s 1}{ }^{(0)}\right\rangle=0 . \tag{3.67}
\end{equation*}
$$

That may not be obvious, but if one returns to the original derivation, the right terms cancel so that one will not end up with the $0 / 0$ problem.

FIXME: performing this derivation outside of class (below), it was found that we do not need the matrix elements of $H^{\prime}$ to be diagonal, but just need

$$
\begin{equation*}
\left\langle\psi_{s \alpha}{ }^{(0)}\right| H^{\prime}\left|\psi_{s \beta}{ }^{(0)}\right\rangle=0, \quad \text { for } \beta \neq \alpha \tag{3.68}
\end{equation*}
$$

That is consistent with problem set III where we did not diagonalize $H^{\prime}$, but just the subset of it associated with the degenerate states. I am unsure now if eq. (3.67) was copied in error or provided in error in class, but it definitely appears to be a more severe requirement than actually needed to deal with perturbation of a state found in a degenerate energy level.

Time independent perturbation with degeneracy. Now we repeat the derivation of the first order perturbation with degenerate states from lecture 4 . We see explicitly how we would get into (divide by zero) trouble if the state we were perturbing had degeneracy. Here I alter the previous derivation to show this explicitly. Like the non-degenerate case, we are covering the time independent perturbation methods from $\S 16.1$ of the text [4].

We start with a known Hamiltonian $H_{0}$, and alter it with the addition of a "small" perturbation

$$
\begin{equation*}
H=H_{0}+\lambda H^{\prime}, \quad \lambda \in[0,1] . \tag{3.69}
\end{equation*}
$$

For the original operator, we assume that a complete set of eigenvectors and eigenkets is known

$$
\begin{equation*}
H_{0}\left|\psi_{s \alpha}{ }^{(0)}\right\rangle=E_{s}{ }^{(0)}\left|\psi_{s \alpha}{ }^{(0)}\right\rangle . \tag{3.70}
\end{equation*}
$$

We seek the perturbed eigensolution

$$
\begin{equation*}
H\left|\psi_{s \alpha}\right\rangle=E_{s \alpha}\left|\psi_{s \alpha}\right\rangle, \tag{3.71}
\end{equation*}
$$

and assumed a perturbative series representation for the energy eigenvalues in the new system

$$
\begin{equation*}
E_{s \alpha}=E_{s}{ }^{(0)}+\lambda E_{s \alpha}{ }^{(1)}+\lambda^{2} E_{s \alpha}{ }^{(2)}+\cdots \tag{3.72}
\end{equation*}
$$

Note that we do not assume that the perturbed energy states, if degenerate in the original system, are still degenerate after perturbation.

Given an assumed representation for the new eigenkets in terms of the known basis

$$
\begin{equation*}
\left|\psi_{s \alpha}\right\rangle=\sum_{n, \beta} c_{n s ; \beta \alpha}\left|\psi_{n \beta}{ }^{(0)}\right\rangle, \tag{3.73}
\end{equation*}
$$

and a perturbative series representation for the probability coefficients

$$
\begin{equation*}
c_{n s ; \beta \alpha}=c_{n s ; \beta \alpha}{ }^{(0)}+\lambda c_{n s ; \beta \alpha}{ }^{(1)}+\lambda^{2} c_{n s ; \beta \alpha}{ }^{(2)}, \tag{3.74}
\end{equation*}
$$

so that

$$
\begin{align*}
\left|\psi_{s \alpha}\right\rangle= & \sum_{n, \beta} c_{n s ; \beta \alpha}{ }^{(0)}\left|\psi_{n \beta}{ }^{(0)}\right\rangle+\lambda \sum_{n, \beta} c_{n s ; \beta \alpha}{ }^{(1)}\left|\psi_{n \beta}{ }^{(0)}\right\rangle \\
& +\lambda^{2} \sum_{n, \beta} c_{n s ; \beta \alpha}{ }^{(2)}\left|\psi_{n \beta}{ }^{(0)}\right\rangle+\cdots \tag{3.75}
\end{align*}
$$

Setting $\lambda=0$ requires

$$
\begin{equation*}
c_{n s ; \beta \alpha}{ }^{(0)}=\delta_{n s ; \beta \alpha}, \tag{3.76}
\end{equation*}
$$

for

$$
\begin{align*}
\left|\psi_{s \alpha}\right\rangle= & \left.\left|\psi_{s \alpha}^{(0)}\right\rangle+\lambda \sum_{n, \beta} c_{n s ; \beta \alpha}{ }^{(1)}\left|\psi_{n \beta^{(0)}}^{(0)}+\lambda^{2} \sum_{n, \beta} c_{n s ; \beta \alpha}^{(2)}\right| \psi_{n \beta^{(0)}}^{(0)}\right\rangle+\cdots \\
= & \left(1+\lambda c_{s s ; \alpha \alpha}^{(1)}+\lambda^{2} c_{s s ; \alpha \alpha}{ }^{(2)}+\cdots\right)\left|\psi_{s \alpha}^{(0)}\right\rangle \\
& +\lambda \sum_{n \beta \neq s \alpha} c_{n s ; \beta \alpha}^{(1)}\left|\psi_{n \beta^{(0)}}^{(0)}\right\rangle \\
& +\lambda^{2} \sum_{n \beta \neq s \alpha} c_{n s ; \beta \alpha}{ }^{(2)}\left|\psi_{n \beta}{ }^{(0)}\right\rangle+\cdots \tag{3.77}
\end{align*}
$$

We rescale our kets

$$
\begin{equation*}
\left|\bar{\psi}_{s \alpha}\right\rangle=\left|\psi_{s \alpha}{ }^{(0)}\right\rangle+\lambda \sum_{n \beta \neq s \alpha} \bar{c}_{n s ; \beta \alpha}{ }^{(1)}\left|\psi_{n \beta}{ }^{(0)}\right\rangle+\lambda^{2} \sum_{n \beta \neq s \alpha} \bar{c}_{n s ; \beta \alpha}{ }^{(2)}\left|\psi_{n \beta}{ }^{(0)}\right\rangle+\cdots \tag{3.78}
\end{equation*}
$$

where

$$
\begin{equation*}
\bar{c}_{n s ; \beta \alpha}{ }^{(j)}=\frac{c_{n s ; \beta \alpha}{ }^{(j)}}{1+\lambda c_{s s ; \alpha \alpha}^{(1)}+\lambda^{2} c_{s s ; \alpha \alpha}^{(2)}+\cdots} \tag{3.79}
\end{equation*}
$$

The normalization of the rescaled kets is then

$$
\begin{equation*}
\left\langle\bar{\psi}_{s \alpha} \mid \bar{\psi}_{s \alpha}\right\rangle=1+\lambda^{2} \sum_{n \beta \neq s \alpha}\left|\bar{c}_{s s}^{(1)}\right|^{2}+\cdots \equiv \frac{1}{Z_{s \alpha}} \tag{3.80}
\end{equation*}
$$

One can then construct a renormalized ket if desired

$$
\begin{equation*}
\left|\bar{\psi}_{s \alpha}\right\rangle_{R}=Z_{s \alpha}^{1 / 2}\left|\bar{\psi}_{s \alpha}\right\rangle \tag{3.81}
\end{equation*}
$$

so that

$$
\begin{equation*}
\left(\left|\bar{\psi}_{s \alpha}\right\rangle_{R}\right)^{\dagger}\left|\bar{\psi}_{s \alpha}\right\rangle_{R}=Z_{s \alpha}\left\langle\bar{\psi}_{s \alpha} \mid \bar{\psi}_{s \alpha}\right\rangle=1 \tag{3.82}
\end{equation*}
$$

We continue by renaming terms in eq. (3.78)

$$
\begin{equation*}
\left|\bar{\psi}_{s \alpha}\right\rangle=\left|\psi_{s \alpha}^{(0)}\right\rangle+\lambda\left|\psi_{s \alpha}^{(1)}\right\rangle+\lambda^{2}\left|\psi_{s \alpha}^{(2)}\right\rangle+\cdots \tag{3.83}
\end{equation*}
$$

where

$$
\begin{equation*}
\left|\psi_{s \alpha}^{(j)}\right\rangle=\sum_{n \beta \neq s \alpha} \bar{c}_{n s ; \beta \alpha}{ }^{(j)}\left|\psi_{n \beta}^{(0)}\right\rangle . \tag{3.84}
\end{equation*}
$$

Now we act on this with the Hamiltonian

$$
\begin{equation*}
H\left|\bar{\psi}_{s \alpha}\right\rangle=E_{s \alpha}\left|\bar{\psi}_{s \alpha}\right\rangle \tag{3.85}
\end{equation*}
$$

or

$$
\begin{equation*}
H\left|\bar{\psi}_{s \alpha}\right\rangle-E_{s \alpha}\left|\bar{\psi}_{s \alpha}\right\rangle=0 . \tag{3.86}
\end{equation*}
$$

Expanding this, we have

$$
\begin{align*}
0= & \left(H_{0}+\lambda H^{\prime}\right)\left(\left|\psi_{s \alpha}{ }^{(0)}\right\rangle+\lambda\left|\psi_{s \alpha}{ }^{(1)}\right\rangle+\lambda^{2}\left|\psi_{s \alpha}{ }^{(2)}\right\rangle+\cdots\right) \\
& -\left(E_{s}{ }^{(0)}+\lambda E_{s \alpha}{ }^{(1)}+\lambda^{2} E_{s \alpha}{ }^{(2)}+\cdots\right) \times  \tag{3.87}\\
& \left(\left|\psi_{s \alpha}{ }^{(0)}\right\rangle+\lambda\left|\psi_{s \alpha}{ }^{(1)}\right\rangle+\lambda^{2}\left|\psi_{s \alpha}{ }^{(2)}\right\rangle+\cdots\right) .
\end{align*}
$$

We want to write this as

$$
\begin{equation*}
|A\rangle+\lambda|B\rangle+\lambda^{2}|C\rangle+\cdots=0 . \tag{3.88}
\end{equation*}
$$

This is

$$
\begin{align*}
0 & =\lambda^{0}\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s \alpha}^{(0)}\right\rangle \\
& +\lambda\left(\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s \alpha}^{(1)}\right\rangle+\left(H^{\prime}-E_{s \alpha}^{(1)}\right)\left|\psi_{s \alpha}^{(0)}\right\rangle\right) \\
& +\lambda^{2}\left(\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s \alpha}^{(2)}\right\rangle+\left(H^{\prime}-E_{s \alpha}^{(1)}\right)\left|\psi_{s \alpha}^{(1)}\right\rangle-E_{s \alpha}^{(2)}\left|\psi_{s \alpha}^{(0)}\right\rangle\right) \tag{3.89}
\end{align*}
$$

So we form

$$
\begin{align*}
& |A\rangle=\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s \alpha}{ }^{(0)}\right\rangle \\
& |B\rangle=\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s \alpha}{ }^{(1)}\right\rangle+\left(H^{\prime}-E_{s \alpha}^{(1)}\right)\left|\psi_{s \alpha}{ }^{(0)}\right\rangle  \tag{3.90}\\
& |C\rangle=\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s \alpha}{ }^{(2)}\right\rangle+\left(H^{\prime}-E_{s \alpha}^{(1)}\right)\left|\psi_{s \alpha}{ }^{(1)}\right\rangle-E_{s \alpha}^{(2)}\left|\psi_{s \alpha}{ }^{(0)}\right\rangle,
\end{align*}
$$

and so forth.
Zeroth order in $\lambda . \quad$ Since $H_{0}\left|\psi_{s \alpha}{ }^{(0)}\right\rangle=E_{s}^{(0)}\left|\psi_{s \alpha}{ }^{(0)}\right\rangle$, this first condition on $|A\rangle$ is not much more than a statement that $0-0=0$.

First order in $\lambda$. How about $|B\rangle=0$ ? For this to be zero we require that both of the following are simultaneously zero

$$
\begin{align*}
\left\langle\psi_{s \alpha}^{(0)} \mid B\right\rangle & =0  \tag{3.91}\\
\left\langle\psi_{m \beta^{(0)}}{ }^{(0)} \mid B\right\rangle & =0, \quad m \beta \neq s \alpha
\end{align*}
$$

This first condition is

$$
\begin{equation*}
\left\langle\psi_{s \alpha}{ }^{(0)}\right|\left(H^{\prime}-E_{s \alpha}^{(1)}\right)\left|\psi_{s \alpha}^{(0)}\right\rangle=0 . \tag{3.92}
\end{equation*}
$$

With

$$
\begin{equation*}
\left\langle\psi_{m \beta}{ }^{(0)}\right| H^{\prime}\left|\psi_{s \alpha}^{(0)}\right\rangle \equiv H_{m s ; \beta \alpha}^{\prime} \tag{3.93}
\end{equation*}
$$

or

$$
\begin{equation*}
H_{s s ; \alpha \alpha}^{\prime}=E_{s \alpha}^{(1)} \tag{3.94}
\end{equation*}
$$

From the second condition we have

$$
\begin{equation*}
0=\left\langle\psi_{m \beta^{(0)}}{ }^{(0)}\right|\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s \alpha}^{(1)}\right\rangle+\left\langle\psi_{m \beta^{(0)}}{ }^{(0)}\right|\left(H^{\prime}-E_{s \alpha}^{(1)}\right)\left|\psi_{s \alpha}^{(0)}\right\rangle . \tag{3.95}
\end{equation*}
$$

Utilizing the Hermitian nature of $H_{0}$ we can act backwards on $\left\langle\psi_{m}{ }^{(0)}\right|$

$$
\begin{equation*}
\left\langle\psi_{m \beta^{(0)}}\right| H_{0}=E_{m}^{(0)}\left\langle\psi_{m \beta}{ }^{(0)}\right| \tag{3.96}
\end{equation*}
$$

We note that $\left\langle\psi_{m \beta}{ }^{(0)} \mid \psi_{s \alpha}{ }^{(0)}\right\rangle=0, m \beta \neq s \alpha$. We can also expand the $\left\langle\psi_{m \beta}{ }^{(0)} \mid \psi_{s \alpha}{ }^{(1)}\right\rangle$, which is

$$
\begin{equation*}
\left.\left\langle\psi_{m \beta^{(0)}}{ }^{(0)} \psi_{s \alpha}^{(1)}\right\rangle=\left\langle\psi_{m \beta^{(0)}}{ }^{(0)}\right| \sum_{n \delta \neq s \alpha} \bar{c}_{n s ; \delta \alpha}^{(1)}\left|\psi_{n \delta}{ }^{(0)}\right\rangle\right) . \tag{3.97}
\end{equation*}
$$

I found that reducing this sum was not obvious until some actual integers were plugged in. Suppose that $s=31$, and $m \beta=22$, then this is

$$
\begin{align*}
\left\langle\psi_{22}{ }^{(0)} \mid \psi_{31}{ }^{(1)}\right\rangle & \left.=\left\langle\psi_{22}{ }^{(0)}\right| \sum_{n \delta \in\{11,12, \cdots, 2} \sum_{1,22,23, \cdots, 32,33, \cdots\}} \bar{c}_{n 3 ; \delta 1}{ }^{(1)}\left|\psi_{n \delta}{ }^{(0)}\right\rangle\right) \\
& =\bar{c}_{23 ; 21}{ }^{(1)}\left\langle\psi_{22}{ }^{(0)} \mid \psi_{22}{ }^{(0)}\right\rangle \\
& =\bar{c}_{23 ; 21}{ }^{(1)} . \tag{3.98}
\end{align*}
$$

Observe that we can also replace the superscript (1) with $(j)$ in the above manipulation without impacting anything else. That and putting back in the abstract indices, we have the general result

$$
\begin{equation*}
\left\langle\psi_{m \beta}{ }^{(0)} \mid \psi_{s \alpha}{ }^{(j)}\right\rangle=\bar{c}_{m s ; \beta \alpha}{ }^{(j)} \tag{3.99}
\end{equation*}
$$

Utilizing this gives us, for $m \beta \neq s \alpha$

$$
\begin{equation*}
0=\left(E_{m}^{(0)}-E_{s}^{(0)}\right) \bar{c}_{m s ; \beta \alpha}{ }^{(1)}+H_{m s ; \beta \alpha}^{\prime} . \tag{3.100}
\end{equation*}
$$

Here we see our first sign of the trouble hinted at in lecture 5. Just because $m \beta \neq s \alpha$ does not mean that $m \neq s$. For example, with $m \beta=11$ and $s \alpha=12$ we would have

$$
\begin{align*}
E_{12}^{(1)} & ={H^{\prime}}_{11 ; 22} \\
\bar{c}_{11 ; 12}{ }^{(1)} & =\frac{H_{11 ; 12}^{\prime}}{E_{1}^{(0)}-E_{1}^{(0)}} . \tag{3.101}
\end{align*}
$$

We have got a divide by zero unless additional restrictions are imposed!
If we return to eq. (3.100), we see that, for the result to be valid, when $m=s$, and there exists degeneracy for the $s$ state, we require for $\beta \neq \alpha$

$$
\begin{equation*}
H_{s s ; \beta \alpha}^{\prime}=0, \tag{3.102}
\end{equation*}
$$

(then eq. (3.100) becomes a $0=0$ equality, and all is still okay)
Summarizing what we learn from our $|B\rangle=0$ conditions we have

$$
\begin{align*}
E_{s \alpha}^{(1)} & ={H^{\prime}}_{s s ; \alpha \alpha} \\
\bar{c}_{m s ; \beta \alpha}^{(1)} & =\frac{{H^{\prime}}_{m s ; \beta \alpha}}{E_{s}^{(0)}-E_{m}^{(0)}}, \quad m \neq s  \tag{3.103}\\
H_{s s ; \beta \alpha}^{\prime} & =0, \quad \beta \alpha \neq 11 .
\end{align*}
$$

Second order in $\lambda$. Doing the same thing for $|C\rangle=0$ we form (or assume)

$$
\begin{equation*}
\left\langle\psi_{s \alpha}^{(0)} \mid C\right\rangle=0 \tag{3.104}
\end{equation*}
$$

and find

$$
\begin{aligned}
0= & \left\langle\psi_{s \alpha}^{(0)} \mid C\right\rangle \\
= & \left\langle\psi_{s \alpha}^{(0)}\right|\left(\left(H_{0}-E_{s}^{(0)}\right)\left|\psi_{s \alpha}^{(2)}\right\rangle+\left(H^{\prime}-E_{s \alpha}^{(1)}\right)\left|\psi_{s \alpha}^{(1)}\right\rangle-E_{s \alpha}^{(2)}\left|\psi_{s \alpha}^{(0)}\right\rangle\right) \\
= & \left(E_{s}^{(0)}-E_{s}^{(0)}\right)\left\langle\psi_{s \alpha}^{(0)} \mid \psi_{s \alpha}^{(2)}\right\rangle+\left\langle\psi_{s \alpha}^{(0)}\right|\left(H^{\prime}-E_{s \alpha}^{(1)}\right)\left|\psi_{s \alpha}^{(1)}\right\rangle \\
& -E_{s \alpha}^{(2)}\left\langle\psi_{s \alpha}^{(0)} \mid \psi_{s \alpha}^{(0)}\right\rangle .
\end{aligned}
$$

We need to know what the $\left\langle\psi_{s \alpha}{ }^{(0)} \mid \psi_{s \alpha}{ }^{(1)}\right\rangle$ is, and find that it is zero

$$
\begin{equation*}
\left\langle\psi_{s \alpha}{ }^{(0)} \mid \psi_{s \alpha}^{(1)}\right\rangle=\left\langle\psi_{s \alpha}{ }^{(0)}\right| \sum_{n \beta \neq s \alpha} \bar{c}_{n s ; \beta \alpha}{ }^{(1)}\left|\psi_{n \beta}{ }^{(0)}\right\rangle=0 \tag{3.106}
\end{equation*}
$$

Utilizing that we have

$$
\begin{align*}
E_{s \alpha}^{(2)} & =\left\langle\psi_{s \alpha}^{(0)}\right| H^{\prime}\left|\psi_{s \alpha}^{(1)}\right\rangle \\
& =\left\langle\psi_{s \alpha}^{(0)}\right| H^{\prime} \sum_{m \beta \neq s \alpha} \bar{c}_{m s}^{(1)}\left|\psi_{m \beta^{(0)}}\right\rangle  \tag{3.107}\\
& =\sum_{m \beta \neq s \alpha} \bar{c}_{m s ; \beta \alpha}^{(1)} H^{\prime}{ }_{s m ; \alpha \beta}
\end{align*}
$$

From eq. (3.103), treating the $m \neq s$ case carefully, we have

$$
\begin{equation*}
E_{s \alpha}^{(2)}=\sum_{\beta \neq \alpha} \bar{c}_{s s ; \beta \alpha}{ }^{(1)}{H^{\prime}}_{s s ; \alpha \beta}+\sum_{m \beta \neq s \alpha, m \neq s} \frac{H_{m s ; \beta \alpha}^{\prime}}{E_{s}^{(0)}-E_{m}^{(0)}} H_{s m ; \alpha \beta}^{\prime} \tag{3.108}
\end{equation*}
$$

Again, only if $H_{s s ; \alpha \beta}=0$ for $\beta \neq \alpha$ do we have a result we can use. If that is the case, the first sum is killed without a divide by zero, leaving

$$
\begin{equation*}
E_{s \alpha}^{(2)}=\sum_{m \beta \neq s \alpha, m \neq s} \frac{\left|H_{m s ; \beta \alpha}^{\prime}\right|^{2}}{E_{s}^{(0)}-E_{m}^{(0)}} \tag{3.109}
\end{equation*}
$$

We can now summarize by forming the first order terms of the perturbed energy and the corresponding kets

$$
\begin{align*}
E_{s \alpha} & =E_{s}^{(0)}+\lambda{H^{\prime}}_{s s ; \alpha \alpha}+\lambda^{2} \sum_{m \neq s, m \beta \neq s \alpha} \frac{\left|H_{m s ; \beta \alpha}^{\prime}\right|^{2}}{E_{s}^{(0)}-E_{m}^{(0)}}+\cdots \\
\left|\bar{\psi}_{s \alpha}\right\rangle & =\left|\psi_{s \alpha}^{(0)}\right\rangle+\lambda \sum_{m \neq s, m \beta \neq s \alpha} \frac{H_{m s ; \beta \alpha}^{\prime}}{E_{s}^{(0)}-E_{m}^{(0)}}\left|\psi_{m \beta}^{(0)}\right\rangle+\cdots  \tag{3.110}\\
H_{s s ; \beta \alpha}^{\prime} & =0, \quad \beta \alpha \neq 11 .
\end{align*}
$$

Notational discrepancy: OOPS. It looks like I used different notation than in class for our matrix elements for the placement of the indices. FIXME: looks like the $c_{s s ; \alpha \alpha^{\prime}}{ }^{(1)}$, for $\alpha \neq \alpha^{\prime}$ coefficients have been lost track of here? Do we have to assume those are zero too? Professor Sipe did not include those in his lecture eq. (3.116), but I do not see the motivation here for dropping them in this derivation.

Diagonalizing the perturbation Hamiltonian. Suppose that we do not have this special zero condition that allows the perturbation treatment to remain valid. What can we do. It turns out that we can make use of the fact that the perturbation Hamiltonian is Hermitian, and diagonalize the matrix

$$
\begin{equation*}
\left\langle\psi_{s \alpha}^{(0)}\right| H^{\prime}\left|\psi_{s \beta}{ }^{(0)}\right\rangle . \tag{3.111}
\end{equation*}
$$

In the example of a two fold degeneracy, this amounts to us choosing not to work with the states

$$
\begin{equation*}
\left|\psi_{s 1}^{(0)}\right\rangle,\left|\psi_{s 2}^{(0)}\right\rangle, \tag{3.112}
\end{equation*}
$$

both some linear combinations of the two

$$
\begin{align*}
& \left|\psi_{s I}^{(0)}\right\rangle=a_{1}\left|\psi_{s 1}^{(0)}\right\rangle+b_{1}\left|\psi_{s 2}^{(0)}\right\rangle  \tag{3.113}\\
& \left|\psi_{s I I}^{(0)}\right\rangle=a_{2}\left|\psi_{s 1}^{(0)}\right\rangle+b_{2}\left|\psi_{s 2}^{(0)}\right\rangle . \tag{3.114}
\end{align*}
$$

In this new basis, once found, we have

$$
\begin{equation*}
\left\langle\psi_{s \alpha}{ }^{(0)}\right| H^{\prime}\left|\psi_{s \beta^{(0)}}{ }^{(0)}\right\rangle=\mathcal{H}_{\alpha} \delta_{\alpha \beta} . \tag{3.115}
\end{equation*}
$$

Utilizing this to fix the previous, one would get if the analysis was repeated correctly

$$
\begin{align*}
E_{s \alpha} & =E_{s}^{(0)}+\lambda H^{\prime}{ }_{s \alpha ; s \alpha}+\lambda^{2} \sum_{m \neq s, \beta} \frac{\left|H_{m \beta ; s \alpha}^{\prime}\right|^{2}}{E_{s}^{(0)}-E_{m}^{(0)}}+\cdots  \tag{3.116}\\
\left|\bar{\psi}_{s \alpha}\right\rangle & =\left|\psi_{s \alpha}^{(0)}\right\rangle+\lambda \sum_{m \neq s, \beta} \frac{H_{m \beta}^{\prime}}{E_{s}^{(0)}-E_{m}^{(0)}}\left|\psi_{m \beta}^{(0)}\right\rangle+\cdots \tag{3.117}
\end{align*}
$$

FIXME: why do we have second order in $\lambda$ terms for the energy when we found those exactly by diagonalization? We found there that the perturbed energy eigenvalues were multivalued with values $E_{s \alpha}=E_{s}^{(0)}+\lambda H^{\prime}{ }_{s \beta ; s \beta}$ for all degeneracy indices $\beta$. Have to repeat the derivation for these more carefully to understand this apparent discrepancy.

We see that a degenerate state can be split by applying perturbation.
FIXME: diagram. $E_{s}^{(0)}$ as one energy level without perturbation, and as two distinct levels with perturbation.

Guess. I will bet that this is the origin of the spectral line splitting, especially given that an atom like hydrogen has degenerate states.

### 3.3 Examples.

## Example 3.1: Stark Shift.

Reading: $\S 16.5$ of [4]. Perturbed the Hamiltonian as follows

$$
\begin{equation*}
H=H_{0}+\lambda H^{\prime} \tag{3.118}
\end{equation*}
$$

where

$$
\begin{equation*}
H^{\prime}=e \mathcal{E}_{z} \hat{Z} \tag{3.119}
\end{equation*}
$$

where $\mathcal{E}_{z}$ is the electric field. To first order the perturbed base state is

$$
\begin{equation*}
\left|\psi_{\alpha}^{(1)}\right\rangle=\left|\psi_{\alpha}^{(0)}\right\rangle+\sum_{\beta \neq \alpha} \frac{\left|\psi_{\beta}^{(0)}\right\rangle\left\langle\psi_{\beta}^{(0)}\right| H^{\prime}\left|\psi_{\alpha}^{(0)}\right\rangle}{E_{\alpha}^{(0)}-E_{\beta}^{(0)}} \tag{3.120}
\end{equation*}
$$

and

$$
\begin{equation*}
E_{\alpha}^{(1)}=\left\langle\psi_{\alpha}^{(0)}\right| H^{\prime}\left|\psi_{\alpha}^{(0)}\right\rangle \tag{3.121}
\end{equation*}
$$

With the default basis $\left\{\left|\psi_{\beta}^{(0)}\right\rangle\right\}$, and $n=2$ we have a 4 fold degeneracy

$$
\begin{align*}
& l, m=0,0 \\
& l, m=1,-1  \tag{3.122}\\
& l, m=1,0 \\
& l, m=1,+1
\end{align*}
$$

but can diagonalize as follows

$$
\left[\begin{array}{ccccc}
\text { nlm } & 200 & 210 & 211 & 21-1  \tag{3.123}\\
200 & 0 & \Delta & 0 & 0 \\
210 & \Delta & 0 & 0 & 0 \\
211 & 0 & 0 & 0 & 0 \\
21-1 & 0 & 0 & 0 & 0
\end{array}\right]
$$

where

$$
\begin{equation*}
\Delta=-3 e \mathcal{E}_{z} a_{0} \tag{3.124}
\end{equation*}
$$

We have a split of energy levels as illustrated in fig. 3.2.


Figure 3.2: Energy level splitting.
Observe the embedded Pauli matrix (I missed the point of that?)

$$
\sigma_{x}=\left[\begin{array}{ll}
0 & 1  \tag{3.125}\\
1 & 0
\end{array}\right]
$$

Proper basis for perturbation (FIXME:check) is then

$$
\begin{equation*}
\left\{\frac{1}{\sqrt{2}}(|2,0,0\rangle \pm|2,1,0\rangle),|2,1, \pm 1\rangle\right\} \tag{3.126}
\end{equation*}
$$

and our result is

$$
\begin{equation*}
\left|\psi_{\alpha, n=2}^{(1)}\right\rangle=\left|\psi_{\alpha}^{(0)}\right\rangle+\sum_{\beta \notin \text { degenerate subspace }} \frac{\left|\psi_{\beta}^{(0)}\right\rangle\left\langle\psi_{\beta}^{(0)}\right| H^{\prime}\left|\psi_{\alpha}^{(0)}\right\rangle}{E_{\alpha}^{(0)}-E_{\beta}^{(0)}} \tag{3.127}
\end{equation*}
$$

## TIME DEPENDENT PERTURBATION.

## 4

### 4.1 REVIEW OF DYNAMICS.

We want to move on to time dependent problems. In general for a time dependent problem, the answer follows provided one has solved for all the perturbed energy eigenvalues. This can be laborious (or not feasible due to infinite sums).

Before doing this, let us review our dynamics as covered in §3 of the text [4].

Schrödinger and Heisenberg pictures. Our operator equation in the Schrödinger picture is the familiar

$$
\begin{equation*}
i \hbar \frac{d}{d t}\left|\psi_{s}(t)\right\rangle=H\left|\psi_{s}(t)\right\rangle \tag{4.1}
\end{equation*}
$$

and most of our operators $X, P, \cdots$ are time independent, but our kets evolve with time

$$
\begin{equation*}
\langle O\rangle(t)=\left\langle\psi_{s}(t)\right| O_{s}\left|\psi_{s}(t)\right\rangle . \tag{4.2}
\end{equation*}
$$

Here $O_{s}$ is the operator in the Schrödinger picture, and is non time dependent. Formally, the time evolution of any state is given by

$$
\begin{equation*}
\left|\psi_{s}(t)\right\rangle e^{-i H t / \hbar}\left|\psi_{s}(0)\right\rangle=U(t, 0)\left|\psi_{s}(0)\right\rangle \tag{4.3}
\end{equation*}
$$

so the expectation of an operator can be written

$$
\begin{equation*}
\langle O\rangle(t)=\left\langle\psi_{s}(0)\right| e^{i H t / \hbar} O_{s} e^{-i H t / \hbar}\left|\psi_{s}(0)\right\rangle . \tag{4.4}
\end{equation*}
$$

With the introduction of the Heisenberg ket

$$
\begin{equation*}
\left|\psi_{H}\right\rangle=\left|\psi_{s}(0)\right\rangle, \tag{4.5}
\end{equation*}
$$

and Heisenberg operators

$$
\begin{equation*}
O_{H}=e^{i H t / \hbar} O_{s} e^{-i H t / \hbar} \tag{4.6}
\end{equation*}
$$

the expectation evolution takes the form

$$
\begin{equation*}
\langle O\rangle(t)=\left\langle\psi_{H}\right| O_{H}\left|\psi_{H}\right\rangle \tag{4.7}
\end{equation*}
$$

Note that because the Hamiltonian commutes with its exponential (it commutes with itself and any power series of itself), the Hamiltonian in the Heisenberg picture is the same as in the Schrödinger picture

$$
\begin{equation*}
H_{H}=e^{i H t / \hbar} H e^{-i H t / \hbar}=H \tag{4.8}
\end{equation*}
$$

Time evolution and the Commutator. Taking the derivative of eq. (4.6) provides us with the time evolution of any operator in the Heisenberg picture

$$
\begin{align*}
i \hbar \frac{d}{d t} O_{H}(t) & =i \hbar \frac{d}{d t}\left(e^{i H t / \hbar} O_{s} e^{-i H t / \hbar}\right) \\
& =i \hbar\left(\frac{i H}{\hbar} e^{i H t / \hbar} O_{s} e^{-i H t / \hbar}+e^{i H t / \hbar} O_{s} e^{\left.-i H t / \hbar \frac{-i H}{\hbar}\right)}\right.  \tag{4.9}\\
& =\left(-H O_{H}+O_{H} H\right)
\end{align*}
$$

We can write this as a commutator

$$
\begin{equation*}
i \hbar \frac{d}{d t} O_{H}(t)=\left[O_{H}, H\right] \tag{4.10}
\end{equation*}
$$

Summarizing the two pictures.
Schrödinger picture Heisenberg picture

$$
\begin{align*}
i \hbar \frac{d}{d t}\left|\psi_{s}(t)\right\rangle=H\left|\psi_{s}(t)\right\rangle & \quad i \hbar \frac{d}{d t} O_{H}(t)=\left[O_{H}, H\right] \\
\left\langle\psi_{s}(t)\right| O_{S}\left|\psi_{s}(t)\right\rangle & =\left\langle\psi_{H}\right| O_{H}\left|\psi_{H}\right\rangle  \tag{4.11}\\
\left|\psi_{s}(0)\right\rangle & =\left|\psi_{H}\right\rangle \\
O_{S} & =O_{H}(0)
\end{align*}
$$

## 4.2 interaction picture.

Recap. Recall our table comparing our two interaction pictures
Schrödinger picture Heisenberg picture

$$
\begin{align*}
i \hbar \frac{d}{d t}\left|\psi_{s}(t)\right\rangle=H\left|\psi_{s}(t)\right\rangle & \quad i \hbar \frac{d}{d t} O_{H}(t)=\left[O_{H}, H\right] \\
\left\langle\psi_{s}(t)\right| O_{S}\left|\psi_{s}(t)\right\rangle & =\left\langle\psi_{H}\right| O_{H}\left|\psi_{H}\right\rangle  \tag{4.12}\\
\left|\psi_{s}(0)\right\rangle & =\left|\psi_{H}\right\rangle \\
O_{S} & =O_{H}(0)
\end{align*}
$$

A motivating example. While fundamental Hamiltonians are independent of time, in a number of common cases, we can form approximate Hamiltonians that are time dependent. One such example is that of Coulomb excitations of an atom, as covered in $\S 18.3$ of the text [4], and shown in fig. 4.1. We consider the interaction of a nucleus with a neutral atom,


Figure 4.1: Coulomb interaction of a nucleus and heavy atom.
heavy enough that it can be considered classically. From the atoms point of view, the effects of the heavy nucleus barrelling by can be described using a time dependent Hamiltonian. For the atom, that interaction Hamiltonian is

$$
\begin{equation*}
H^{\prime}=\sum_{i} \frac{Z e q_{i}}{\left|\mathbf{r}_{N}(t)-\mathbf{R}_{i}\right|} \tag{4.13}
\end{equation*}
$$

Here and $\mathbf{r}_{N}$ is the position vector for the heavy nucleus, and $\mathbf{R}_{i}$ is the position to each charge within the atom, where $i$ ranges over all the internal charges, positive and negative, within the atom.

Placing the origin close to the atom, we can write this interaction Hamiltonian as

$$
\begin{equation*}
H^{\prime}(t)=\sum_{j} \frac{Z e q_{i}}{\left|\mathbf{r}_{N}(t)\right|}+\left.\sum_{i} Z e q_{i} \mathbf{R}_{i} \cdot\left(\frac{\partial}{\partial \mathbf{r}} \frac{1}{\left|\mathbf{r}_{N}(t)-\mathbf{r}\right|}\right)\right|_{\mathbf{r}=0} \tag{4.14}
\end{equation*}
$$

The first term vanishes because the total charge in our neutral atom is zero. This leaves us with

$$
\begin{align*}
H^{\prime}(t) & =-\left.\sum_{i} q_{i} \mathbf{R}_{i} \cdot\left(-\frac{\partial}{\partial \mathbf{r}} \frac{Z e}{\left|\mathbf{r}_{N}(t)-\mathbf{r}\right|}\right)\right|_{\mathbf{r}=0}  \tag{4.15}\\
& =-\sum_{i} q_{i} \mathbf{R}_{i} \cdot \mathbf{E}(t)
\end{align*}
$$

where $\mathbf{E}(t)$ is the electric field at the origin due to the nucleus. Introducing a dipole moment operator for the atom

$$
\begin{equation*}
\boldsymbol{\mu}=\sum_{i} q_{i} \mathbf{R}_{i} \tag{4.16}
\end{equation*}
$$

the interaction takes the form

$$
\begin{equation*}
H^{\prime}(t)=-\boldsymbol{\mu} \cdot \mathbf{E}(t) \tag{4.17}
\end{equation*}
$$

Here we have a quantum mechanical operator, and a classical field taken together. This sort of dipole interaction also occurs when we treat a atom placed into an electromagnetic field, treated classically as depicted in fig. 4.2. In the figure, we can use the dipole interaction, provided $\lambda \gg a$,


Figure 4.2: atom in a field.
where $a$ is the "width" of the atom. Because it is great for examples, we will see this dipole interaction a lot.

The interaction picture. Having talked about both the Schrödinger and Heisenberg pictures, we can now move on to describe a hybrid, one where our Hamiltonian has been split into static and time dependent parts

$$
\begin{equation*}
H(t)=H_{0}+H^{\prime}(t) \tag{4.18}
\end{equation*}
$$

We will formulate an approach for dealing with problems of this sort called the interaction picture.

This is also covered in $\S 3.3$ of the text, albeit in a much harder to understand fashion (the text appears to try to not pull the result from a magic hat, but the steps to get to the end result are messy). It would probably have been nicer to see it this way instead.

In the Schrödinger picture our dynamics have the form

$$
\begin{equation*}
i \hbar \frac{d}{d t}\left|\psi_{s}(t)\right\rangle=H\left|\psi_{s}(t)\right\rangle \tag{4.19}
\end{equation*}
$$

How about the Heisenberg picture? We look for a solution

$$
\begin{equation*}
\left|\psi_{s}(t)\right\rangle=U\left(t, t_{0}\right)\left|\psi_{s}\left(t_{0}\right)\right\rangle \tag{4.20}
\end{equation*}
$$

We want to find this operator that evolves the state from the state as some initial time $t_{0}$, to the arbitrary later state found at time $t$. Plugging in we have

$$
\begin{equation*}
i \hbar \frac{d}{d t} U\left(t, t_{0}\right)\left|\psi_{s}\left(t_{0}\right)\right\rangle=H(t) U\left(t, t_{0}\right)\left|\psi_{s}\left(t_{0}\right)\right\rangle \tag{4.21}
\end{equation*}
$$

This has to hold for all $\left|\psi_{s}\left(t_{0}\right)\right\rangle$, and we can equivalently seek a solution of the operator equation

$$
\begin{equation*}
i \hbar \frac{d}{d t} U\left(t, t_{0}\right)=H(t) U\left(t, t_{0}\right) \tag{4.22}
\end{equation*}
$$

where

$$
\begin{equation*}
U\left(t_{0}, t_{0}\right)=I, \tag{4.23}
\end{equation*}
$$

the identity for the Hilbert space. Suppose that $H(t)$ was independent of time. We could find that

$$
\begin{equation*}
U\left(t, t_{0}\right)=e^{-i H\left(t-t_{0}\right) / \hbar} . \tag{4.24}
\end{equation*}
$$

If $H(t)$ depends on time could you guess that

$$
\begin{equation*}
U\left(t, t_{0}\right)=e^{-\frac{i}{\hbar} \int_{t_{0}}^{t} H(\tau) d \tau} \tag{4.25}
\end{equation*}
$$

holds? No. This may be true when $H(t)$ is a number, but when it is an operator, the Hamiltonian does not necessarily commute with itself at different times

$$
\begin{equation*}
\left[H\left(t^{\prime}\right), H\left(t^{\prime \prime}\right)\right] \neq 0 \tag{4.26}
\end{equation*}
$$

So this is wrong in general. As an aside, for numbers, eq. (4.25) can be verified easily. We have

$$
\begin{align*}
i \hbar\left(e^{-\frac{i}{\hbar} \int_{t_{0}}^{t} H(\tau) d \tau}\right)^{\prime} & =i \hbar\left(-\frac{i}{\hbar}\right)\left(\int_{t_{0}}^{t} H(\tau) d \tau\right)^{\prime} e^{-\frac{i}{\hbar} \int_{t_{0}}^{t} H(\tau) d \tau} \\
& =\left(H(t) \frac{d t}{d t}-H\left(t_{0}\right) \frac{d t_{0}}{d t}\right) e^{-\frac{i}{\hbar} \int_{t_{0}}^{t} H(\tau) d \tau}  \tag{4.27}\\
& =H(t) U\left(t, t_{0}\right)
\end{align*}
$$

Expectations. Suppose that we do find $U\left(t, t_{0}\right)$. Then our expectation takes the form

$$
\begin{equation*}
\left\langle\psi_{s}(t)\right| O_{s}\left|\psi_{s}(t)\right\rangle=\left\langle\psi_{s}\left(t_{0}\right)\right| U^{\dagger}\left(t, t_{0}\right) O_{s} U\left(t, t_{0}\right)\left|\psi_{s}\left(t_{0}\right)\right\rangle \tag{4.28}
\end{equation*}
$$

Put

$$
\begin{equation*}
\left|\psi_{H}\right\rangle=\left|\psi_{s}\left(t_{0}\right)\right\rangle \tag{4.29}
\end{equation*}
$$

and form

$$
\begin{equation*}
O_{H}=U^{\dagger}\left(t, t_{0}\right) O_{s} U\left(t, t_{0}\right) \tag{4.30}
\end{equation*}
$$

so that our expectation has the familiar representations

$$
\begin{equation*}
\left\langle\psi_{s}(t)\right| O_{s}\left|\psi_{s}(t)\right\rangle=\left\langle\psi_{H}\right| O_{H}\left|\psi_{H}\right\rangle \tag{4.31}
\end{equation*}
$$

New strategy. Interaction picture. Let us define

$$
\begin{equation*}
U_{I}\left(t, t_{0}\right)=e^{\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} U\left(t, t_{0}\right) \tag{4.32}
\end{equation*}
$$

or

$$
\begin{equation*}
U\left(t, t_{0}\right)=e^{-\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} U_{I}\left(t, t_{0}\right) \tag{4.33}
\end{equation*}
$$

Let us see how this works. We have

$$
\begin{align*}
i \hbar \frac{d U_{I}}{d t} & =i \hbar \frac{d}{d t}\left(e^{\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} U\left(t, t_{0}\right)\right) \\
& =-H_{0} U\left(t, t_{0}\right)+e^{\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)}\left(i \hbar \frac{d}{d t} U\left(t, t_{0}\right)\right) \\
& =-H_{0} U\left(t, t_{0}\right)+e^{\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)}\left(\left(H+H^{\prime}(t)\right) U\left(t, t_{0}\right)\right)  \tag{4.34}\\
& =e^{\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} H^{\prime}(t) U\left(t, t_{0}\right) \\
& =e^{\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} H^{\prime}(t) e^{-\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} U_{I}\left(t, t_{0}\right) .
\end{align*}
$$

Define

$$
\begin{equation*}
\bar{H}^{\prime}(t)=e^{\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} H^{\prime}(t) e^{-\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)}, \tag{4.35}
\end{equation*}
$$

so that our operator equation takes the form

$$
\begin{equation*}
i \hbar \frac{d}{d t} U_{I}\left(t, t_{0}\right)=\bar{H}^{\prime}(t) U_{I}\left(t, t_{0}\right) . \tag{4.36}
\end{equation*}
$$

Note that we also have the required identity at the initial time

$$
\begin{equation*}
U_{I}\left(t_{0}, t_{0}\right)=I . \tag{4.37}
\end{equation*}
$$

Without requiring us to actually find $U\left(t, t_{0}\right)$ all of the dynamics of the time dependent interaction are now embedded in our operator equation for $\bar{H}^{\prime}(t)$, with all of the simple interaction related to the non time dependent portions of the Hamiltonian left separate.

Connection with the Schrödinger picture. In the Schrödinger picture we have

$$
\begin{align*}
\left|\psi_{s}(t)\right\rangle & =U\left(t, t_{0}\right)\left|\psi_{s}\left(t_{0}\right)\right\rangle \\
& =e^{-\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} U_{I}\left(t, t_{0}\right)\left|\psi_{s}\left(t_{0}\right)\right\rangle . \tag{4.38}
\end{align*}
$$

With a definition of the interaction picture ket as

$$
\begin{equation*}
\left|\psi_{I}\right\rangle=U_{I}\left(t, t_{0}\right)\left|\psi_{s}\left(t_{0}\right)\right\rangle=U_{I}\left(t, t_{0}\right)\left|\psi_{H}\right\rangle \tag{4.39}
\end{equation*}
$$

the Schrödinger picture is then related to the interaction picture by

$$
\begin{equation*}
\left|\psi_{s}(t)\right\rangle=e^{-\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)}\left|\psi_{I}\right\rangle . \tag{4.40}
\end{equation*}
$$

Also, by multiplying eq. (4.36) by our Schrödinger ket, we remove the last vestiges of $U_{I}$ and $U$ from the dynamical equation for our time dependent interaction

$$
\begin{equation*}
i \hbar \frac{d}{d t}\left|\psi_{I}\right\rangle=\bar{H}^{\prime}(t)\left|\psi_{I}\right\rangle . \tag{4.41}
\end{equation*}
$$

Interaction picture expectation. Inverting eq. (4.40), we can form an operator expectation, and relate it the interaction and Schrödinger pictures

$$
\begin{equation*}
\left\langle\psi_{s}(t)\right| O_{s}\left|\psi_{s}(t)\right\rangle=\left\langle\psi_{I}\right| e^{\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} O_{s} e^{-\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)}\left|\psi_{I}\right\rangle \tag{4.42}
\end{equation*}
$$

With a definition

$$
\begin{equation*}
O_{I}=e^{\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} O_{s} e^{-\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} \tag{4.43}
\end{equation*}
$$

we have

$$
\begin{equation*}
\left\langle\psi_{s}(t)\right| O_{s}\left|\psi_{s}(t)\right\rangle=\left\langle\psi_{I}\right| O_{I}\left|\psi_{I}\right\rangle \tag{4.44}
\end{equation*}
$$

As before, the time evolution of our interaction picture operator, can be found by taking derivatives of eq. (4.43), for which we find

$$
\begin{equation*}
i \hbar \frac{d O_{I}(t)}{d t}=\left[O_{I}(t), H_{0}\right] \tag{4.45}
\end{equation*}
$$

Summarizing the interaction picture. Given

$$
\begin{equation*}
H(t)=H_{0}+H^{\prime}(t) \tag{4.46}
\end{equation*}
$$

and initial time states

$$
\begin{equation*}
\left|\psi_{I}\left(t_{0}\right)\right\rangle=\left|\psi_{s}\left(t_{0}\right)\right\rangle=\left|\psi_{H}\right\rangle \tag{4.47}
\end{equation*}
$$

we have

$$
\begin{equation*}
\left\langle\psi_{s}(t)\right| O_{s}\left|\psi_{s}(t)\right\rangle=\left\langle\psi_{I}\right| O_{I}\left|\psi_{I}\right\rangle \tag{4.48}
\end{equation*}
$$

where

$$
\begin{equation*}
\left|\psi_{I}\right\rangle=U_{I}\left(t, t_{0}\right)\left|\psi_{s}\left(t_{0}\right)\right\rangle \tag{4.49}
\end{equation*}
$$

and

$$
\begin{equation*}
i \hbar \frac{d}{d t}\left|\psi_{I}\right\rangle=\bar{H}^{\prime}(t)\left|\psi_{I}\right\rangle \tag{4.50}
\end{equation*}
$$

or

$$
\begin{align*}
i \hbar \frac{d}{d t} U_{I}\left(t, t_{0}\right) & =\bar{H}^{\prime}(t) U_{I}\left(t, t_{0}\right)  \tag{4.51}\\
U_{I}\left(t_{0}, t_{0}\right) & =I
\end{align*}
$$

Our interaction picture Hamiltonian is

$$
\begin{equation*}
\bar{H}^{\prime}(t)=e^{\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} H^{\prime}(t) e^{-\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)}, \tag{4.52}
\end{equation*}
$$

and for Schrödinger operators, independent of time, we have the dynamical equation

$$
\begin{equation*}
i \hbar \frac{d O_{I}(t)}{d t}=\left[O_{I}(t), H_{0}\right] . \tag{4.53}
\end{equation*}
$$

## 4.3 justifying the taylor expansion above (not class notes).

Multivariable Taylor series. As outlined in §2.8 (8.10) of [9], we want to derive the multi-variable Taylor expansion for a scalar valued function of some number of variables

$$
\begin{equation*}
f(\mathbf{u})=f\left(u^{1}, u^{2}, \cdots\right), \tag{4.54}
\end{equation*}
$$

consider the displacement operation applied to the vector argument

$$
\begin{equation*}
f(\mathbf{a}+\mathbf{x})=\left.f(\mathbf{a}+t \mathbf{x})\right|_{t=1} . \tag{4.55}
\end{equation*}
$$

We can Taylor expand a single variable function without any trouble, so introduce

$$
\begin{equation*}
g(t)=f(\mathbf{a}+t \mathbf{x}) \tag{4.56}
\end{equation*}
$$

where

$$
\begin{equation*}
g(1)=f(\mathbf{a}+\mathbf{x}) . \tag{4.57}
\end{equation*}
$$

We have

$$
\begin{equation*}
g(t)=g(0)+\left.t \frac{\partial g}{\partial t}\right|_{t=0}+\left.\frac{t^{2}}{2!} \frac{\partial g}{\partial t}\right|_{t=0}+\cdots, \tag{4.58}
\end{equation*}
$$

so that

$$
\begin{equation*}
g(1)=g(0)++\left.\frac{\partial g}{\partial t}\right|_{t=0}+\left.\frac{1}{2!} \frac{\partial g}{\partial t}\right|_{t=0}+\cdots \tag{4.59}
\end{equation*}
$$

The multivariable Taylor series now becomes a plain old application of the chain rule, where we have to evaluate

$$
\begin{align*}
\frac{d g}{d t} & =\frac{d}{d t} f\left(a^{1}+t x^{1}, a^{2}+t x^{2}, \cdots\right) \\
& =\sum_{i} \frac{\partial}{\partial\left(a^{i}+t x^{i}\right)} f(\mathbf{a}+t \mathbf{x}) \frac{\partial a^{i}+t x^{i}}{\partial t} \tag{4.60}
\end{align*}
$$

so that

$$
\begin{equation*}
\left.\frac{d g}{d t}\right|_{t=0}=\sum_{i} x^{i}\left(\left.\frac{\partial f}{\partial x^{i}}\right|_{x^{i}=a^{i}}\right) . \tag{4.61}
\end{equation*}
$$

Assuming an Euclidean space we can write this in the notationally more pleasant fashion using a gradient operator for the space

$$
\begin{equation*}
\left.\frac{d g}{d t}\right|_{t=0}=\left.\mathbf{x} \cdot \nabla_{\mathbf{u}} f(\mathbf{u})\right|_{\mathbf{u}=\mathbf{a}} . \tag{4.62}
\end{equation*}
$$

To handle the higher order terms, we repeat the chain rule application, yielding for example

$$
\begin{align*}
\left.\frac{d^{2} f(\mathbf{a}+t \mathbf{x})}{d t^{2}}\right|_{t=0} & =\left.\frac{d}{d t} \sum_{i} x^{i} \frac{\partial f(\mathbf{a}+t \mathbf{x})}{\partial\left(a^{i}+t x^{i}\right)}\right|_{t=0} \\
& =\left.\sum_{i} x^{i} \frac{\partial}{\partial\left(a^{i}+t x^{i}\right)} \frac{d f(\mathbf{a}+t \mathbf{x})}{d t}\right|_{t=0}  \tag{4.63}\\
& =\left.\left(\mathbf{x} \cdot \nabla_{\mathbf{u}}\right)^{2} f(\mathbf{u})\right|_{\mathbf{u}=\mathbf{a}} .
\end{align*}
$$

Thus the Taylor series associated with a vector displacement takes the tidy form

$$
\begin{equation*}
f(\mathbf{a}+\mathbf{x})=\left.\sum_{k=0}^{\infty} \frac{1}{k!}\left(\mathbf{x} \cdot \nabla_{\mathbf{u}}\right)^{k} f(\mathbf{u})\right|_{\mathbf{u}=\mathbf{a}} . \tag{4.64}
\end{equation*}
$$

Even more fancy, we can form the operator equation

$$
\begin{equation*}
f(\mathbf{a}+\mathbf{x})=\left.e^{\mathbf{x} \cdot \nabla_{\mathbf{u}}} f(\mathbf{u})\right|_{\mathbf{u}=\mathbf{a}} . \tag{4.65}
\end{equation*}
$$

Here a dummy variable $\mathbf{u}$ has been retained as an instruction not to differentiate the $\mathbf{x}$ part of the directional derivative in any repeated applications of the $\mathbf{x} \cdot \boldsymbol{\nabla}$ operator.

That notational kludge can be removed by swapping $\mathbf{a}$ and $\mathbf{x}$

$$
\begin{equation*}
f(\mathbf{a}+\mathbf{x})=\sum_{k=0}^{\infty} \frac{1}{k!}(\mathbf{a} \cdot \boldsymbol{\nabla})^{k} f(\mathbf{x})=e^{\mathbf{a} \cdot \boldsymbol{\nabla}} f(\mathbf{x}), \tag{4.66}
\end{equation*}
$$

where $\boldsymbol{\nabla}=\boldsymbol{\nabla}_{\mathbf{x}}=\left(\partial / \partial x^{1}, \partial / \partial x^{2}, \ldots\right)$.
Having derived this (or for those with lesser degrees of amnesia, recall it), we can see that eq. (4.14) was a direct application of this, retaining no second order or higher terms.

Our expression used in the interaction Hamiltonian discussion was

$$
\begin{equation*}
\frac{1}{|\mathbf{r}-\mathbf{R}|} \approx \frac{1}{|\mathbf{r}|}+\left.\mathbf{R} \cdot\left(\frac{\partial}{\partial \mathbf{R}} \frac{1}{|\mathbf{r}-\mathbf{R}|}\right)\right|_{\mathbf{R}=0} \tag{4.67}
\end{equation*}
$$

which we can see has the same structure as above with some variable substitutions. Evaluating it we have

$$
\begin{align*}
\frac{\partial}{\partial \mathbf{R}} \frac{1}{|\mathbf{r}-\mathbf{R}|} & =\mathbf{e}_{i} \frac{\partial}{\partial R^{i}}\left(\left(x^{j}-R^{j}\right)^{2}\right)^{-1 / 2} \\
& =\mathbf{e}_{i}\left(-\frac{1}{2}\right) 2\left(x^{j}-R^{j}\right) \frac{\partial\left(x^{j}-R^{j}\right)}{\partial R^{i}} \frac{1}{|\mathbf{r}-\mathbf{R}|^{3}}  \tag{4.68}\\
& =\frac{\mathbf{r}-\mathbf{R}}{|\mathbf{r}-\mathbf{R}|^{3}},
\end{align*}
$$

and at $\mathbf{R}=0$ we have

$$
\begin{equation*}
\frac{1}{|\mathbf{r}-\mathbf{R}|} \approx \frac{1}{|\mathbf{r}|}+\mathbf{R} \cdot \frac{\mathbf{r}}{|\mathbf{r}|^{3}} \tag{4.69}
\end{equation*}
$$

We see in this direction derivative produces the classical electric Coulomb field expression for an electrostatic distribution, once we take the $\mathbf{r} /|\mathbf{r}|^{3}$ and multiply it with the $-Z e$ factor.

With algebra. A different way to justify the expansion of eq. (4.14) is to consider a Clifford algebra factorization (following notation from [7]) of the absolute vector difference, where $\mathbf{R}$ is considered small.

$$
\begin{align*}
|\mathbf{r}-\mathbf{R}| & =\sqrt{(\mathbf{r}-\mathbf{R})(\mathbf{r}-\mathbf{R})} \\
& =\sqrt{\left\langle\mathbf{r}\left(1-\frac{1}{\mathbf{r}} \mathbf{R}\right)\left(1-\mathbf{R} \frac{1}{\mathbf{r}}\right) \mathbf{r}\right\rangle} \\
& =\sqrt{\left\langle\mathbf{r}^{2}\left(1-\frac{1}{\mathbf{r}} \mathbf{R}\right)\left(1-\mathbf{R} \frac{1}{\mathbf{r}}\right)\right\rangle}  \tag{4.70}\\
& =|\mathbf{r}| \sqrt{1-2 \frac{1}{\mathbf{r}} \cdot \mathbf{R}+\left\langle\frac{1}{\mathbf{r}} \mathbf{R} \frac{1}{\mathbf{r}}\right\rangle} \\
& =|\mathbf{r}| \sqrt{1-2 \frac{1}{\mathbf{r}} \cdot \mathbf{R}+\frac{\mathbf{R}^{2}}{\mathbf{r}^{2}}}
\end{align*}
$$

Neglecting the $\mathbf{R}^{2}$ term, we can then Taylor series expand this scalar expression

$$
\begin{equation*}
\frac{1}{|\mathbf{r}-\mathbf{R}|} \approx \frac{1}{|\mathbf{r}|}\left(1+\frac{1}{\mathbf{r}} \cdot \mathbf{R}\right)=\frac{1}{|\mathbf{r}|}+\frac{\hat{\mathbf{r}}}{\mathbf{r}^{2}} \cdot \mathbf{R}=\frac{1}{|\mathbf{r}|}+\frac{\mathbf{r}}{|\mathbf{r}|^{3}} \cdot \mathbf{R} . \tag{4.71}
\end{equation*}
$$

Observe this is what was found with the multivariable Taylor series expansion too.

## 4.4 recap: interaction picture.

We will use the interaction picture to examine time dependent perturbations. We wrote our Schrödinger ket in terms of the interaction ket

$$
\begin{equation*}
|\psi\rangle=e^{-i H_{0}\left(t-t_{0}\right) / \hbar}\left|\psi_{I}(t)\right\rangle, \tag{4.72}
\end{equation*}
$$

where

$$
\begin{equation*}
\left|\psi_{I}\right\rangle=U_{I}\left(t, t_{0}\right)\left|\psi_{I}\left(t_{0}\right)\right\rangle . \tag{4.73}
\end{equation*}
$$

Our dynamics is given by the operator equation

$$
\begin{equation*}
i \hbar \frac{d}{d t} U_{I}\left(t, t_{0}\right)=\bar{H}^{\prime}(t) U_{I}\left(t, t_{0}\right), \tag{4.74}
\end{equation*}
$$

where

$$
\begin{equation*}
\bar{H}^{\prime}(t)=e^{\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} H^{\prime}(t) e^{-\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} . \tag{4.75}
\end{equation*}
$$

We can formally solve eq. (4.74) by writing

$$
\begin{equation*}
U_{I}\left(t, t_{0}\right)=I-\frac{i}{\hbar} \int_{t_{0}}^{t} d t^{\prime} \bar{H}^{\prime}\left(t^{\prime}\right) U_{I}\left(t^{\prime}, t_{0}\right) \tag{4.76}
\end{equation*}
$$

This is easy enough to verify by direct differentiation

$$
\begin{align*}
i \hbar \frac{d}{d t} U_{I} & =\left(\int_{t_{0}}^{t} d t^{\prime} \bar{H}^{\prime}\left(t^{\prime}\right) U_{I}\left(t^{\prime}, t_{0}\right)\right)^{\prime} \\
& =\bar{H}^{\prime}(t) U_{I}\left(t, t_{0}\right) \frac{d t}{d t}-\bar{H}^{\prime}(t) U_{I}\left(t, t_{0}\right) \frac{d t_{0}}{d t}  \tag{4.77}\\
& =\bar{H}^{\prime}(t) U_{I}\left(t, t_{0}\right) .
\end{align*}
$$

This is a bit of a chicken and an egg expression, since it is cyclic with a dependency on unknown $U_{I}\left(t^{\prime}, t_{0}\right)$ factors. We start with an initial estimate of the operator to be determined, and iterate. This can seem like an odd thing to do, but one can find books on just this integral kernel iteration method (like the nice little Dover book [16] that has sat on my (Peeter's) shelf all lonely so many years). Suppose for $t$ near $t_{0}$, try

$$
\begin{equation*}
U_{I}\left(t, t_{0}\right) \approx I-\frac{i}{\hbar} \int_{t_{0}}^{t} d t^{\prime} \bar{H}^{\prime}\left(t^{\prime}\right) \tag{4.78}
\end{equation*}
$$

A second order iteration is now possible

$$
\begin{align*}
U_{I}\left(t, t_{0}\right) & \approx I-\frac{i}{\hbar} \int_{t_{0}}^{t} d t^{\prime} \bar{H}^{\prime}\left(t^{\prime}\right)\left(I-\frac{i}{\hbar} \int_{t_{0}}^{t^{\prime}} d t^{\prime \prime} \bar{H}^{\prime}\left(t^{\prime \prime}\right) .\right) \\
& =I-\frac{i}{\hbar} \int_{t_{0}}^{t} d t^{\prime} \bar{H}^{\prime}\left(t^{\prime}\right)+\left(\frac{-i}{\hbar}\right)^{2} \int_{t_{0}}^{t} d t^{\prime} \bar{H}^{\prime}\left(t^{\prime}\right) \int_{t_{0}}^{t^{\prime}} d t^{\prime \prime} \bar{H}^{\prime}\left(t^{\prime \prime}\right) \tag{4.79}
\end{align*}
$$

It is possible to continue this iteration, and this approach is considered in some detail in $\S 3.3$ of the text [4], and is apparently also the basis for Feynman diagrams.

### 4.5 TIME DEPENDENT PERTURBATION THEORY.

As covered in $\S 17$ of the text, we will split the interaction into time independent and time dependent terms

$$
\begin{equation*}
H(t)=H_{0}+H^{\prime}(t) \tag{4.80}
\end{equation*}
$$

and work in the interaction picture with

$$
\begin{equation*}
\left|\psi_{I}(t)\right\rangle=\sum_{n} \tilde{c}_{n}(t)\left|\psi_{n}^{(0)}\right\rangle \tag{4.81}
\end{equation*}
$$

Our Schrödinger ket is then

$$
\begin{align*}
|\psi(t)\rangle & =e^{-i H_{0}^{(0)}\left(t-t_{0}\right) / \hbar}\left|\psi_{I}\left(t_{0}\right)\right\rangle \\
& =\sum_{n} \tilde{c}_{n}(t) e^{-i E_{n}^{(0)}\left(t-t_{0}\right) / \hbar}\left|\psi_{n}^{(0)}\right\rangle . \tag{4.82}
\end{align*}
$$

With a definition

$$
\begin{equation*}
c_{n}(t)=\tilde{c}_{n}(t) e^{i E_{n} t_{0} / \hbar} \tag{4.83}
\end{equation*}
$$

(where we leave off the zero superscript for the unperturbed state), our time evolved ket becomes

$$
\begin{equation*}
|\psi(t)\rangle=\sum_{n} c_{n}(t) e^{-i E_{n} t / \hbar}\left|\psi_{n}^{(0)}\right\rangle \tag{4.84}
\end{equation*}
$$

We can now plug eq. (4.81) into our evolution equation

$$
\begin{align*}
i \hbar \frac{d}{d t}\left|\psi_{I}(t)\right\rangle & =\bar{H}^{\prime}(t)\left|\psi_{I}(t)\right\rangle  \tag{4.85}\\
& =e^{\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} H^{\prime}(t) e^{-\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)}\left|\psi_{I}(t)\right\rangle
\end{align*}
$$

which gives us

$$
\begin{equation*}
i \hbar \sum_{p} \frac{\partial}{\partial t} \tilde{c}_{p}(t)\left|\psi_{p}^{(0)}\right\rangle=e^{\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} H^{\prime}(t) e^{-\frac{i}{\hbar} H_{0}\left(t-t_{0}\right)} \sum_{n} \tilde{c}_{n}(t)\left|\psi_{n}^{(0)}\right\rangle . \tag{4.86}
\end{equation*}
$$

We can apply the bra $\left\langle\psi_{m}^{(0)}\right|$ to this equation, yielding

$$
\begin{equation*}
i \hbar \frac{\partial}{\partial t} \tilde{c}_{m}(t)=\sum_{n} \tilde{c}_{n}(t) e^{\frac{i}{\hbar} E_{m}\left(t-t_{0}\right)}\left\langle\psi_{m}^{(0)}\right| H^{\prime}(t)\left|\psi_{n}^{(0)}\right\rangle e^{-\frac{i}{\hbar} E_{n}\left(t-t_{0}\right)} \tag{4.87}
\end{equation*}
$$

With

$$
\begin{align*}
\omega_{m} & =\frac{E_{m}}{\hbar} \\
\omega_{m n} & =\omega_{m}-\omega_{n}  \tag{4.88}\\
H_{m n}^{\prime}(t) & =\left\langle\psi_{m}^{(0)}\right| H^{\prime}(t)\left|\psi_{n}^{(0)}\right\rangle,
\end{align*}
$$

this is

$$
\begin{equation*}
i \hbar \frac{\partial \tilde{c}_{m}(t)}{\partial t}=\sum_{n} \tilde{c}_{n}(t) e^{i \omega_{m n}\left(t-t_{0}\right)} H_{m n}^{\prime}(t) \tag{4.89}
\end{equation*}
$$

Inverting eq. (4.83) and plugging in

$$
\begin{equation*}
\tilde{c}_{n}(t)=c_{n}(t) e^{-i \omega_{n} t_{0}} \tag{4.90}
\end{equation*}
$$

yields

$$
\begin{equation*}
i \hbar \frac{\partial c_{m}(t)}{\partial t} e^{-i \omega_{m} t_{0}}=\sum_{n} c_{n}(t) e^{-i \omega_{n} t_{0}} e^{i \omega_{m n} t} e^{-i\left(\omega_{m}-\omega_{n}\right) t_{0}} H_{m n}^{\prime}(t) \tag{4.91}
\end{equation*}
$$

from which we can cancel the exponentials on both sides yielding

$$
\begin{equation*}
i \hbar \frac{\partial c_{m}(t)}{\partial t}=\sum_{n} c_{n}(t) e^{i \omega_{m n} t} H_{m n}^{\prime}(t) \tag{4.92}
\end{equation*}
$$

We are now left with all of our time dependence nicely separated out, with the coefficients $c_{n}(t)$ encoding all the non-oscillatory time evolution information

$$
\begin{align*}
H & =H_{0}+H^{\prime}(t) \\
|\psi(t)\rangle & =\sum_{n} c_{n}(t) e^{-i \omega_{n} t}\left|\psi_{n}^{(0)}\right\rangle  \tag{4.93}\\
i \hbar \dot{c}_{m} & =\sum_{n} H_{m n}^{\prime}(t) e^{i \omega_{m n} t} c_{n}(t) .
\end{align*}
$$

### 4.6 PERTURBATION EXPANSION.

We now introduce our $\lambda$ parametrization

$$
\begin{equation*}
H^{\prime}(t) \rightarrow \lambda H^{\prime}(t), \tag{4.94}
\end{equation*}
$$

and hope for convergence, or at least something that at least has well defined asymptotic behavior. We have

$$
\begin{equation*}
i \hbar \dot{c}_{m}=\lambda \sum_{n} H_{m n}^{\prime}(t) e^{i \omega_{m n} t} c_{n}(t) \tag{4.95}
\end{equation*}
$$

and try

$$
\begin{equation*}
c_{m}(t)=c_{m}^{(0)}(t)+\lambda c_{m}^{(1)}(t)+\lambda^{2} c_{m}^{(2)}(t)+\cdots \tag{4.96}
\end{equation*}
$$

Plugging in, we have

$$
\begin{equation*}
i \hbar \sum_{k} \lambda^{k} \dot{c}_{m}^{(k)}(t)=\sum_{n, p} H_{m n}^{\prime}(t) e^{i \omega_{m n} t} \lambda^{p+1} c_{n}^{(p)}(t) \tag{4.97}
\end{equation*}
$$

As before, for equality, we treat this as an equation for each $\lambda^{k}$. Expanding explicitly for the first few powers, gives us

$$
\begin{align*}
0 & =\lambda^{0}\left(i \hbar \dot{c}_{m}^{(0)}(t)-0\right) \\
& +\lambda^{1}\left(i \hbar \dot{c}_{m}^{(1)}(t)-\sum_{n} H_{m n}^{\prime}(t) e^{i \omega_{m n} t} c_{n}^{(0)}(t)\right) \\
& +\lambda^{2}\left(i \hbar \dot{c}_{m}^{(2)}(t)-\sum_{n} H_{m n}^{\prime}(t) e^{i \omega_{m n} t} c_{n}^{(1)}(t)\right) \tag{4.98}
\end{align*}
$$

Suppose we have a set of energy levels as depicted in fig. 4.3. With $c_{n}^{(i)}=0$ before the perturbation for all $i \geq 1, n$ and $c_{m}^{(0)}=\delta_{m s}$, we can proceed iteratively, solving each equation, starting with

$$
\begin{equation*}
i \hbar \dot{c}_{m}^{(1)}=H_{m s}^{\prime}(t) e^{i \omega_{m s} t} . \tag{4.99}
\end{equation*}
$$

## Example 4.1: Slow nucleus passing an atom.

Let

$$
\begin{equation*}
H^{\prime}(t)=-\mu \cdot \mathbf{E}(t) \tag{4.100}
\end{equation*}
$$



Figure 4.3: Perturbation around energy level s.
where

$$
\begin{equation*}
H_{m s}^{\prime}=-\mu_{m s} \cdot \mathbf{E}(t), \tag{4.101}
\end{equation*}
$$

and

$$
\begin{equation*}
\boldsymbol{\mu}_{m s}=\left\langle\psi_{m}^{(0)}\right| \boldsymbol{\mu}\left|\psi_{s}^{(0)}\right\rangle . \tag{4.102}
\end{equation*}
$$

Using our previous nucleus passing an atom example, as depicted in fig. 4.4.


Figure 4.4: Slow nucleus passing an atom.
We have

$$
\begin{equation*}
\boldsymbol{\mu}=\sum_{i} q_{i} \mathbf{R}_{i}, \tag{4.103}
\end{equation*}
$$

the dipole moment for each of the charges in the atom. We will have fields as depicted in fig. 4.5.


Figure 4.5: Fields for nucleus atom example.

## Example 4.2: Electromagnetic wave pulse interacting with an atom.

Consider a EM wave pulse, perhaps Gaussian, of the form depicted in fig. 4.6.


Figure 4.6: Atom interacting with an EM pulse.

$$
\begin{equation*}
E_{y}(t)=e^{-t^{2} / T^{2}} \cos \left(\omega_{0} t\right) \tag{4.104}
\end{equation*}
$$

As we learned very early, perhaps sitting on our mother's knee, we can solve the differential equation eq. (4.99) for the first order perturbation, by direct integration

$$
\begin{equation*}
c_{m}^{(1)}(t)=\frac{1}{i \hbar} \int_{-\infty}^{t} H_{m s}^{\prime}\left(t^{\prime}\right) e^{i \omega_{m s} t^{\prime}} d t^{\prime} \tag{4.105}
\end{equation*}
$$

Here the perturbation is assumed equal to zero at $-\infty$. Suppose our electric field is specified in terms of a Fourier transform

$$
\begin{equation*}
\mathbf{E}(t)=\int_{-\infty}^{\infty} \frac{d \omega}{2 \pi} \mathbf{E}(\omega) e^{-i \omega t} \tag{4.106}
\end{equation*}
$$

so

$$
\begin{equation*}
c_{m}^{(1)}(t)=\frac{\boldsymbol{\mu}_{m s}}{2 \pi i \hbar} \cdot \int_{-\infty}^{\infty} \int_{-\infty}^{t} \mathbf{E}(\omega) e^{i\left(\omega_{m s}-\omega\right) t^{\prime}} d t^{\prime} d \omega \tag{4.107}
\end{equation*}
$$

From this, "after the perturbation", as $t \rightarrow \infty$ we find

$$
\begin{align*}
c_{m}^{(1)}(\infty) & =\frac{\boldsymbol{\mu}_{m s}}{2 \pi i \hbar} \cdot \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \mathbf{E}(\omega) e^{i\left(\omega_{m s}-\omega\right) t^{\prime}} d t^{\prime} d \omega \\
& =\frac{\boldsymbol{\mu}_{m s}}{i \hbar} \cdot \int_{-\infty}^{\infty} \mathbf{E}(\omega) \delta\left(\omega_{m s}-\omega\right) d \omega \tag{4.108}
\end{align*}
$$

since we identify

$$
\begin{equation*}
\frac{1}{2 \pi} \int_{-\infty}^{\infty} e^{i\left(\omega_{m s}-\omega\right) t^{\prime}} d t^{\prime} \equiv \delta\left(\omega_{m s}-\omega\right) \tag{4.109}
\end{equation*}
$$

Thus the steady state first order perturbation coefficient is

$$
\begin{equation*}
c_{m}^{(1)}(\infty)=\frac{\boldsymbol{\mu}_{m s}}{i \hbar} \cdot \mathbf{E}\left(\omega_{m s}\right) \tag{4.110}
\end{equation*}
$$

Frequency symmetry for the Fourier spectrum of a real field. We will look further at this next week, but we first require an intermediate result from transform theory. Because our field is real, we have

$$
\begin{equation*}
\mathbf{E}^{*}(t)=\mathbf{E}(t) \tag{4.111}
\end{equation*}
$$

so

$$
\begin{align*}
\mathbf{E}^{*}(t) & =\int \frac{d \omega}{2 \pi} \mathbf{E}^{*}(\omega) e^{i \omega t}  \tag{4.112}\\
& =\int \frac{d \omega}{2 \pi} \mathbf{E}^{*}(-\omega) e^{-i \omega t}
\end{align*}
$$

and thus

$$
\begin{equation*}
\mathbf{E}(\omega)=\mathbf{E}^{*}(-\omega), \tag{4.113}
\end{equation*}
$$

and

$$
\begin{equation*}
|\mathbf{E}(\omega)|^{2}=|\mathbf{E}(-\omega)|^{2} \tag{4.114}
\end{equation*}
$$

We will see shortly what the point of this aside is.

### 4.7 TIME DEPENDENT PERTURBATION.

We would gotten as far as calculating

$$
\begin{equation*}
c_{m}^{(1)}(\infty)=\frac{1}{i \hbar} \boldsymbol{\mu}_{m s} \cdot \mathbf{E}\left(\omega_{m s}\right) \tag{4.115}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathbf{E}(t)=\int \frac{d \omega}{2 \pi} \mathbf{E}(\omega) e^{-i \omega t} \tag{4.116}
\end{equation*}
$$

and

$$
\begin{equation*}
\omega_{m s}=\frac{E_{m}-E_{s}}{\hbar} \tag{4.117}
\end{equation*}
$$

Graphically, these frequencies are illustrated in fig. 4.7. The probability


Figure 4.7: Positive and negative frequencies.
for a transition from $m$ to $s$ is therefore

$$
\begin{equation*}
\rho_{m \rightarrow s}=\left|c_{m}^{(1)}(\infty)\right|^{2}=\frac{1^{2}}{\hbar}\left|\boldsymbol{\mu}_{m s} \cdot \mathbf{E}\left(\omega_{m s}\right)\right|^{2} \tag{4.118}
\end{equation*}
$$

Recall that because the electric field is real we had

$$
\begin{equation*}
|\mathbf{E}(\omega)|^{2}=|\mathbf{E}(-\omega)|^{2} \tag{4.119}
\end{equation*}
$$



Figure 4.8: Gaussian wave packet.

Suppose that we have a wave pulse, where our field magnitude is perhaps of the form

$$
\begin{equation*}
E(t)=e^{-t^{2} / T^{2}} \cos \left(\omega_{0} t\right) \tag{4.120}
\end{equation*}
$$

as illustrated with $\omega=10, T=1$ in fig. 4.8. We expect this to have a two lobe Fourier spectrum, with the lobes centered at $\omega= \pm 10$, and width proportional to $1 / T$.

For reference, as calculated using qmTwoL8figures.nb, this Fourier transform is

$$
\begin{equation*}
E(\omega)=\frac{e^{-\frac{1}{4} T^{2}\left(\omega_{0}+\omega\right)^{2}}}{2 \sqrt{\frac{2}{T^{2}}}}+\frac{e^{\omega_{0} T^{2} \omega-\frac{1}{4} T^{2}\left(\omega_{0}+\omega\right)^{2}}}{2 \sqrt{\frac{2}{T^{2}}}} \tag{4.121}
\end{equation*}
$$

This is illustrated, again for $\omega_{0}=10$, and $T=1$, in fig. 4.9, where we see


Figure 4.9: Fourier transform of Gaussian wave packet.
the expected Gaussian result, since the Fourier transform of a Gaussian is a Gaussian.

### 4.8 Sudden perturbations.

Given our wave equation

$$
\begin{equation*}
i \hbar \frac{d}{d t}|\psi(t)\rangle=H(t)|\psi(t)\rangle, \tag{4.122}
\end{equation*}
$$

and a sudden perturbation in the Hamiltonian, as illustrated in fig. 4.10. Consider $H_{0}$ and $H_{F}$ fixed, and decrease $\Delta t \rightarrow 0$. We can formally inte-


Figure 4.10: Sudden step Hamiltonian.
grate eq. (4.122)

$$
\begin{equation*}
\frac{d}{d t}|\psi(t)\rangle=\frac{1}{i \hbar} H(t)|\psi(t)\rangle . \tag{4.123}
\end{equation*}
$$

For

$$
\begin{equation*}
|\psi(t)\rangle-\left|\psi\left(t_{0}\right)\right\rangle=\frac{1}{i \hbar} \int_{t_{0}}^{t} H\left(t^{\prime}\right)\left|\psi\left(t^{\prime}\right)\right\rangle d t^{\prime} \tag{4.124}
\end{equation*}
$$

While this is an exact solution, it is also not terribly useful since we do not know $|\psi(t)\rangle$. However, we can select the small interval $\Delta t$, and write

$$
\begin{equation*}
|\psi(\Delta t / 2)\rangle=|\psi(-\Delta t / 2)\rangle+\frac{1}{i \hbar} \int_{t_{0}}^{t} H\left(t^{\prime}\right)\left|\psi\left(t^{\prime}\right)\right\rangle d t^{\prime} . \tag{4.125}
\end{equation*}
$$

Note that we could use the integral kernel iteration technique here and substitute $\left|\psi\left(t^{\prime}\right)\right\rangle=|\psi(-\Delta t / 2)\rangle$ and then develop this, to generate a power series with $(\Delta t / 2)^{k}$ dependence. However, we note that eq. (4.125) is still an exact relation, and if $\Delta t \rightarrow 0$, with the integration limits narrowing (provided $H\left(t^{\prime}\right)$ is well behaved) we are left with just

$$
\begin{equation*}
|\psi(\Delta t / 2)\rangle=|\psi(-\Delta t / 2)\rangle \tag{4.126}
\end{equation*}
$$

or

$$
\begin{equation*}
\left|\psi_{\text {after }}\right\rangle=\left|\psi_{\text {before }}\right\rangle \tag{4.127}
\end{equation*}
$$

provided that we change the Hamiltonian fast enough. On the surface there appears to be no consequences, but there are some very serious ones!

## Example 4.3: Harmonic oscillator.

Consider our harmonic oscillator Hamiltonian, with

$$
\begin{align*}
H_{0} & =\frac{P^{2}}{2 m}+\frac{1}{2} m \omega_{0}^{2} X^{2}  \tag{4.128}\\
H_{F} & =\frac{P^{2}}{2 m}+\frac{1}{2} m \omega_{F}^{2} X^{2}
\end{align*}
$$

Here $\omega_{0} \rightarrow \omega_{F}$ continuously, but very quickly. In effect, we have tightened the spring constant. Note that there are cases in linear optics when you can actually do exactly that.

Imagine that $\left|\psi_{\text {before }}\right\rangle$ is in the ground state of the harmonic oscillator as in fig. 4.11,


Figure 4.11: Harmonic oscillator sudden Hamiltonian perturbation.
and we suddenly change the Hamiltonian with potential $V_{0} \rightarrow V_{F}$ (weakening the "spring"). Professor Sipe gives us a graphical demo of this, by impersonating a constrained wavefunction with his arms, doing weak chicken-flapping of them. Now with the potential weakened, he wiggles and flaps his arms with more freedom and somewhat chaotically. His "wave function" arms are now bouncing around in the new limiting potential (initially over doing it and then bouncing back).

We had in this case the exact relation

$$
\begin{equation*}
H_{0}\left|\psi_{0}^{(0)}\right\rangle=\frac{1}{2} \hbar \omega_{0}\left|\psi_{0}^{(0)}\right\rangle \tag{4.129}
\end{equation*}
$$

but we also have

$$
\begin{equation*}
\left|\psi_{\text {after }}\right\rangle=\left|\psi_{\text {before }}\right\rangle=\left|\psi_{0}^{(0)}\right\rangle \tag{4.130}
\end{equation*}
$$

and

$$
\begin{equation*}
H_{F}\left|\psi_{n}^{(f)}\right\rangle=\frac{1}{2} \hbar \omega_{F}\left(n+\frac{1}{2}\right)\left|\psi_{n}^{(f)}\right\rangle \tag{4.131}
\end{equation*}
$$

So

$$
\begin{align*}
\left|\psi_{\mathrm{after}}\right\rangle & =\left|\psi_{0}^{(0)}\right\rangle \\
& =\sum_{n}\left|\psi_{n}^{(f)}\right\rangle\left\langle\psi_{n}^{(f)} \mid \psi_{0}^{(0)}\right\rangle  \tag{4.132}\\
& =\sum_{n} c_{n}\left|\psi_{n}^{(f)}\right\rangle
\end{align*}
$$

and at later times

$$
\begin{align*}
\left|\psi(t)^{(f)}\right\rangle & =\left|\psi_{0}^{(0)}\right\rangle \\
& =\sum_{n} c_{n} e^{i \omega_{n}^{(f)} t}\left|\psi_{n}^{(f)}\right\rangle \tag{4.133}
\end{align*}
$$

whereas

$$
\begin{equation*}
\left|\psi(t)^{(o)}\right\rangle=e^{i \omega_{0}^{(0)} t}\left|\psi_{0}^{(0)}\right\rangle \tag{4.134}
\end{equation*}
$$

So, while the wave functions may be exactly the same after such a sudden change in Hamiltonian, the dynamics of the situation change for all future times, since we now have a wavefunction that has a different set of components in the basis for the new Hamiltonian. In particular, the evolution of the wave function is now significantly more complex.

FIXME: plot an example of this.

### 4.9 ADIABATIC PERTURBATIONS.

This is treated in $\S 17.5 .2$ of the text [4].
I wondered what Adiabatic meant in this context. The usage in class sounds like it was just "really slow and gradual", yet this has a definition [5] "Of, relating to, or being a reversible thermodynamic process that occurs without gain or loss of heat and without a change in entropy". Wikipedia [17] appears to confirm that the QM meaning of this term is just "slow" changing.

This is the reverse case, and we now vary the Hamiltonian $H(t)$ very slowly.

$$
\begin{equation*}
\frac{d}{d t}|\psi(t)\rangle=\frac{1}{i \hbar} H(t)|\psi(t)\rangle . \tag{4.135}
\end{equation*}
$$

We first consider only non-degenerate states, and at $t=0$ write

$$
\begin{equation*}
H(0)=H_{0}, \tag{4.136}
\end{equation*}
$$

and

$$
\begin{equation*}
H_{0}\left|\psi_{s}^{(0)}\right\rangle=E_{s}^{(0)}\left|\psi_{s}^{(0)}\right\rangle \tag{4.137}
\end{equation*}
$$

Imagine that at each time $t$ we can find the "instantaneous" energy eigenstates

$$
\begin{equation*}
H(t)\left|\hat{\psi}_{s}(t)\right\rangle=E_{s}(t)\left|\hat{\psi}_{s}(t)\right\rangle \tag{4.138}
\end{equation*}
$$

These states do not satisfy Schrödinger's equation, but are simply solutions to the eigen problem. Our standard strategy in perturbation is based on analysis of

$$
\begin{equation*}
|\psi(t)\rangle=\sum_{n} c_{n}(t) e^{-i \omega_{n}^{(0)} t}\left|\psi_{n}^{(0)}\right\rangle \tag{4.139}
\end{equation*}
$$

Here instead

$$
\begin{equation*}
|\psi(t)\rangle=\sum_{n} b_{n}(t)\left|\hat{\psi}_{n}(t)\right\rangle \tag{4.140}
\end{equation*}
$$

we will expand, not using our initial basis, but instead using the instantaneous kets. Plugging into Schrödinger's equation we have

$$
\begin{align*}
H(t)|\psi(t)\rangle & =H(t) \sum_{n} b_{n}(t)\left|\hat{\psi}_{n}(t)\right\rangle \\
& =\sum_{n} b_{n}(t) E_{n}(t)\left|\hat{\psi}_{n}(t)\right\rangle . \tag{4.141}
\end{align*}
$$

This was complicated before with matrix elements all over the place. Now it is easy, however, the time derivative becomes harder. Doing that we find

$$
\begin{align*}
i \hbar \frac{d}{d t}|\psi(t)\rangle & =i \hbar \frac{d}{d t} \sum_{n} b_{n}(t)\left|\hat{\psi}_{n}(t)\right\rangle \\
& =i \hbar \sum_{n} \frac{d b_{n}(t)}{d t}\left|\hat{\psi}_{n}(t)\right\rangle+\sum_{n} b_{n}(t) \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle  \tag{4.142}\\
& =\sum_{n} b_{n}(t) E_{n}(t)\left|\hat{\psi}_{n}(t)\right\rangle
\end{align*}
$$

We bra $\left\langle\hat{\psi}_{m}(t)\right|$ into this

$$
\begin{align*}
& i \hbar \sum_{n} \frac{d b_{n}(t)}{d t}\left\langle\hat{\psi}_{m}(t) \mid \hat{\psi}_{n}(t)\right\rangle+\sum_{n} b_{n}(t)\left\langle\hat{\psi}_{m}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle  \tag{4.143}\\
& \quad=\sum_{n} b_{n}(t) E_{n}(t)\left\langle\hat{\psi}_{m}(t) \mid \hat{\psi}_{n}(t)\right\rangle
\end{align*}
$$

and find

$$
\begin{equation*}
i \hbar \frac{d b_{m}(t)}{d t}+\sum_{n} b_{n}(t)\left\langle\hat{\psi}_{m}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle=b_{m}(t) E_{m}(t) \tag{4.144}
\end{equation*}
$$

If the Hamiltonian is changed very very slowly in time, we can imagine that $\left|\hat{\psi}_{n}(t)\right\rangle^{\prime}$ is also changing very very slowly, but we are not quite there yet. Let us first split our sum of bra and ket products

$$
\begin{equation*}
\sum_{n} b_{n}(t)\left\langle\hat{\psi}_{m}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle \tag{4.145}
\end{equation*}
$$

into $n \neq m$ and $n=m$ terms. Looking at just the $n=m$ term

$$
\begin{equation*}
\left\langle\hat{\psi}_{m}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{m}(t)\right\rangle \tag{4.146}
\end{equation*}
$$

We note

$$
\begin{align*}
0 & =\frac{d}{d t}\left\langle\hat{\psi}_{m}(t) \mid \hat{\psi}_{m}(t)\right\rangle \\
& =\left(\frac{d}{d t}\left\langle\hat{\psi}_{m}(t)\right|\right)\left|\hat{\psi}_{m}(t)\right\rangle+\left\langle\hat{\psi}_{m}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{m}(t)\right\rangle \tag{4.147}
\end{align*}
$$

Something plus its complex conjugate equals 0

$$
\begin{equation*}
a+i b+(a+i b)^{*}=2 a=0 \Longrightarrow a=0 \tag{4.148}
\end{equation*}
$$

so $\left\langle\hat{\psi}_{m}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{m}(t)\right\rangle$ must be purely imaginary. We write

$$
\begin{equation*}
\left\langle\hat{\psi}_{m}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{m}(t)\right\rangle=-i \Gamma_{s}(t) \tag{4.149}
\end{equation*}
$$

where $\Gamma_{s}$ is real.

### 4.10 adiabatic perturbation theory (cont.)

We were working through Adiabatic time dependent perturbation (as also covered in §17.5.2 of the text [4].)

Utilizing an expansion

$$
\begin{align*}
|\psi(t)\rangle & =\sum_{n} c_{n}(t) e^{-i \omega_{n}^{(0)} t}\left|\psi_{n}^{(0)}\right\rangle  \tag{4.150}\\
& =\sum_{n} b_{n}(t)\left|\hat{\psi}_{n}(t)\right\rangle
\end{align*}
$$

where

$$
\begin{equation*}
H(t)\left|\hat{\psi}_{s}(t)\right\rangle=E_{s}(t)\left|\hat{\psi}_{s}(t)\right\rangle \tag{4.151}
\end{equation*}
$$

and found

$$
\begin{equation*}
\frac{d b_{s}(t)}{d t}=-i\left(\omega_{s}(t)-\Gamma_{s}(t)\right) b_{s}(t)-\sum_{n \neq s} b_{n}(t)\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle \tag{4.152}
\end{equation*}
$$

where

$$
\begin{equation*}
\Gamma_{s}(t)=i\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{s}(t)\right\rangle \tag{4.153}
\end{equation*}
$$

Look for a solution of the form

$$
\begin{align*}
b_{s}(t) & =\bar{b}_{s}(t) e^{-i \int_{0}^{t} d t^{\prime}\left(\omega_{s}\left(t^{\prime}\right)-\Gamma_{s}\left(t^{\prime}\right)\right)} \\
& =\bar{b}_{s}(t) e^{-i \gamma_{s}(t)} \tag{4.154}
\end{align*}
$$

where

$$
\begin{equation*}
\gamma_{s}(t)=\int_{0}^{t} d t^{\prime}\left(\omega_{s}\left(t^{\prime}\right)-\Gamma_{s}\left(t^{\prime}\right)\right) \tag{4.155}
\end{equation*}
$$

Taking derivatives of $\bar{b}_{s}$ and after a bit of manipulation we find that things conveniently cancel

$$
\begin{align*}
\frac{d \bar{b}_{s}(t)}{d t} & =\frac{d}{d t}\left(b_{s}(t) e^{i \gamma_{s}(t)}\right) \\
& =\frac{d b_{s}(t)}{d t} e^{i \gamma_{s}(t)}+b_{s}(t) \frac{d}{d t} e^{i \gamma_{s}(t)}  \tag{4.156}\\
& =\frac{d b_{s}(t)}{d t} e^{i \gamma_{s}(t)}+b_{s}(t) i\left(\omega_{s}(t)-\Gamma_{s}(t)\right) e^{i \gamma_{s}(t)}
\end{align*}
$$

We find

$$
\begin{align*}
\frac{d \bar{b}_{s}(t)}{d t} e^{-i \gamma_{s}(t)}= & \frac{d b_{s}(t)}{d t}+i b_{s}(t)\left(\omega_{s}(t)-\Gamma_{s}(t)\right) \\
= & \frac{i b_{s}(t)\left(\omega_{s}(t)-\Gamma_{s}(t)\right)}{}-i\left(\omega_{s}(t)=\Gamma_{s}(t)\right) b_{s}(t)  \tag{4.157}\\
& -\sum_{n \neq s} b_{n}(t)\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle,
\end{align*}
$$

so

$$
\begin{align*}
\frac{d \bar{b}_{s}(t)}{d t} & =-\sum_{n \neq s} b_{n}(t) e^{i \gamma_{s}(t)}\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle  \tag{4.158}\\
& =-\sum_{n \neq s} \bar{b}_{n}(t) e^{i\left(\gamma_{s}(t)-\gamma_{n}(t)\right)}\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle .
\end{align*}
$$

With a last bit of notation

$$
\begin{equation*}
\left.\gamma_{s n}(t)=\gamma_{s}(t)-\gamma_{n}(t)\right), \tag{4.159}
\end{equation*}
$$

the problem is reduced to one involving only the sums over the $n \neq s$ terms, and where all the dependence on $\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{s}(t)\right\rangle$ has been nicely isolated in a phase term

$$
\begin{equation*}
\frac{d \bar{b}_{s}(t)}{d t}=-\sum_{n \neq s} \bar{b}_{n}(t) e^{i \gamma_{s n}(t)}\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle . \tag{4.160}
\end{equation*}
$$

Looking for an approximate solution.
Try: An approximate solution

$$
\begin{equation*}
\bar{b}_{n}(t)=\delta_{n m} \tag{4.161}
\end{equation*}
$$

For $s=m$ this is okay, since we have $\frac{d \delta_{n s}}{d t}=0$ which is consistent with

$$
\begin{equation*}
\sum_{n \neq s} \delta_{n s}(\cdots)=0 \tag{4.162}
\end{equation*}
$$

However, for $s \neq m$ we get

$$
\begin{align*}
\frac{d \bar{b}_{s}(t)}{d t} & =-\sum_{n \neq s} \delta_{n m} e^{i \gamma_{s n}(t)}\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle  \tag{4.163}\\
& =-e^{i \gamma_{s m}(t)}\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{m}(t)\right\rangle .
\end{align*}
$$

But

$$
\begin{equation*}
\gamma_{s m}(t)=\int_{0}^{t} d t^{\prime}\left(\frac{1}{\hbar}\left(E_{s}\left(t^{\prime}\right)-E_{m}\left(t^{\prime}\right)\right)-\Gamma_{s}\left(t^{\prime}\right)+\Gamma_{m}\left(t^{\prime}\right)\right) . \tag{4.164}
\end{equation*}
$$

FIXME: I think we argued in class that the $\Gamma$ contributions are negligible. Why was that?

Now, are energy levels will have variation with time, as illustrated in fig. 4.12. Perhaps unrealistically, suppose that our energy levels have


Figure 4.12: Energy level variation with time.
some "typical" energy difference $\Delta E$, so that

$$
\begin{equation*}
\gamma_{s m}(t) \approx \frac{\Delta E}{\hbar} t \equiv \frac{t}{\tau}, \tag{4.165}
\end{equation*}
$$

or

$$
\begin{equation*}
\tau=\frac{\hbar}{\Delta E} \tag{4.166}
\end{equation*}
$$

Suppose that $\tau$ is much less than a typical time $T$ over which instantaneous quantities (wavefunctions and brakets) change. After a large time $T$

$$
\begin{equation*}
e^{i \gamma_{s m}(t)} \approx e^{i T / \tau} \tag{4.167}
\end{equation*}
$$

so we have our phase term whipping around really fast, as illustrated in fig. 4.13. So, while $\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{m}(t)\right\rangle$ is moving really slow, but our phase


Figure 4.13: Phase whipping around.
space portion is changing really fast. The key to the approximate solution is factoring out this quickly changing phase term.

Note. $\quad \Gamma_{s}(t)$ is called the "Berry" phase [18], whereas the $E_{s}\left(t^{\prime}\right) / \hbar$ part is called the geometric phase, and can be shown to have a geometric interpretation. To proceed we can introduce $\lambda$ terms, perhaps

$$
\begin{equation*}
\bar{b}_{s}(t)=\delta_{m s}+\lambda \bar{b}_{s}^{(1)}(t)+\cdots \tag{4.168}
\end{equation*}
$$

and

$$
\begin{equation*}
-\sum_{n \neq s} e^{i \gamma_{s n}(t)} \lambda(\cdots) \tag{4.169}
\end{equation*}
$$

This $\lambda$ approximation and a similar Taylor series expansion in time have been explored further in $E$.

Degeneracy. Suppose we have some branching of energy levels that were initially degenerate, as illustrated in fig. 4.14. We have a necessity
$E_{31}(t), E_{32}(t)$
$E_{21}(t)$
$E_{11}(t), E_{12}(t), E_{13}(t)$
$E_{0}(t)$

Figure 4.14: Degenerate energy level splitting.
to choose states properly so there is a continuous evolution in the instantaneous eigenvalues as $H(t)$ changes.

### 4.11 EXAMPLes.

## Example 4.4: Adiabatic perturbation theory.

Utilizing instantaneous eigenstates

$$
\begin{equation*}
|\psi(t)\rangle=\sum_{\alpha} b_{\alpha}(t)\left|\hat{\psi}_{\alpha}(t)\right\rangle \tag{4.170}
\end{equation*}
$$

where

$$
\begin{equation*}
H(t)\left|\hat{\psi}_{\alpha}(t)\right\rangle=E_{\alpha}(t)\left|\hat{\psi}_{\alpha}(t)\right\rangle . \tag{4.171}
\end{equation*}
$$

## We found

$$
\begin{equation*}
b_{\alpha}(t)=\bar{b}_{\alpha}(t) e^{-\frac{i}{\hbar} \int_{0}^{t}\left(E_{\alpha}\left(t^{\prime}\right)-\hbar \Gamma_{\alpha}\left(t^{\prime}\right)\right) d t^{\prime}}, \tag{4.172}
\end{equation*}
$$

where

$$
\begin{equation*}
\Gamma_{\alpha}=i\left\langle\hat{\psi}_{\alpha}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{\alpha}(t)\right\rangle \tag{4.173}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{d}{d t} \bar{b}_{\alpha}(t)=-\sum_{\beta \neq \alpha} \bar{b}_{\beta}(t) e^{-\frac{i}{\hbar} \int_{0}^{t}\left(E_{\beta \alpha}\left(t^{\prime}\right)-\hbar \Gamma_{\beta \alpha}\left(t^{\prime}\right)\right) d t^{\prime}}\left\langle\hat{\psi}_{\alpha}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle \tag{4.174}
\end{equation*}
$$

Suppose we start in a subspace

$$
\begin{equation*}
\operatorname{span}\left\{\frac{1}{\sqrt{2}}(|2,0,0\rangle \pm|2,1,0\rangle),|2,1, \pm 1\rangle\right\} \tag{4.175}
\end{equation*}
$$

Now expand the bra derivative kets

$$
\begin{gather*}
\left\langle\hat{\psi}_{\alpha}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle=\left(\left\langle\psi_{\alpha}^{(0)}\right|+\sum_{\gamma} \frac{\left\langle\psi_{\gamma}^{(0)}\right| H^{\prime}\left|\psi_{\alpha}^{(0)}\right\rangle\left\langle\psi_{\gamma}^{(0)}\right|}{E_{\alpha}^{(0)}-E_{\gamma}^{(0)}}\right) \times \\
\frac{d}{d t}\left(\left|\psi_{\beta}^{(0)}\right\rangle+\sum_{\gamma^{\prime}} \frac{\left|\psi_{\gamma^{\prime}}^{(0)}\right\rangle\left\langle\psi_{\gamma^{\prime}}^{(0)}\right| H^{\prime}\left|\psi_{\beta}^{(0)}\right\rangle}{E_{\beta}^{(0)}-E_{\gamma^{\prime}}^{(0)}}\right) \tag{4.176}
\end{gather*}
$$

To first order we can drop the quadratic terms in $\gamma, \gamma^{\prime}$ leaving

$$
\begin{align*}
\left\langle\hat{\psi}_{\alpha}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle & \sim \sum_{\gamma^{\prime}}\left\langle\psi_{\alpha}^{(0)} \mid \psi_{\gamma^{\prime}}^{(0)}\right\rangle \frac{\left\langle\psi_{\gamma^{\prime}}^{(0)}\right| \frac{d H^{\prime}(t)}{d t}\left|\psi_{\beta}^{(0)}\right\rangle}{E_{\beta}^{(0)}-E_{\gamma^{\prime}}^{(0)}} \\
& =\frac{\left\langle\psi_{\alpha}^{(0)}\right| \frac{d H^{\prime}(t)}{d t}\left|\psi_{\beta}^{(0)}\right\rangle}{E_{\beta}^{(0)}-E_{\alpha}^{(0)}}, \tag{4.177}
\end{align*}
$$

SO

$$
\begin{equation*}
\frac{d}{d t} \bar{b}_{\alpha}(t)=-\sum_{\beta \neq \alpha} \bar{b}_{\beta}(t) e^{-\frac{i}{\hbar} \int_{0}^{t}\left(E_{\beta \alpha}\left(t^{\prime}\right)-\hbar \Gamma_{\beta \alpha}\left(t^{\prime}\right)\right) d t^{\prime}} \frac{\left\langle\psi_{\alpha}^{(0)}\right| \frac{d H^{\prime}(t)}{d t}\left|\psi_{\beta}^{(0)}\right\rangle}{E_{\beta}^{(0)}-E_{\alpha}^{(0)}} \tag{4.178}
\end{equation*}
$$

A different way to this end result. A result of this form is also derived in [2] §20.1, but with a different approach. There he takes derivatives of

$$
\begin{gather*}
H(t)\left|\hat{\psi}_{\beta}(t)\right\rangle=E_{\beta}(t)\left|\hat{\psi}_{\beta}(t)\right\rangle, \\
\frac{d H(t)}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle+H(t) \frac{d}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle=\frac{d E_{\beta}(t)}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle+E_{\beta}(t) \frac{d}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle . \tag{4.180}
\end{gather*}
$$

Bra'ing $\left\langle\hat{\psi}_{\alpha}(t)\right|$ into this we have, for $\alpha \neq \beta$, the LHS is

$$
\begin{align*}
& \left\langle\hat{\psi}_{\alpha}(t)\right| \frac{d H(t)}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle+\left\langle\hat{\psi}_{\alpha}(t)\right| H(t) \frac{d}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle  \tag{4.181}\\
& \quad=\left\langle\hat{\psi}_{\alpha}(t)\right| \frac{d H(t)}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle+E_{\alpha}(t)\left\langle\hat{\psi}_{\alpha}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle,
\end{align*}
$$

and for the RHS

$$
\begin{equation*}
\left\langle\hat{\psi}_{\alpha}(t)\right| \frac{d E_{\beta}(t)}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle+\left\langle\hat{\psi}_{\alpha}(t)\right| E_{\beta}(t) \frac{d}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle, \tag{4.182}
\end{equation*}
$$

so

$$
\begin{equation*}
\left\langle\hat{\psi}_{\alpha}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle=\frac{\left\langle\hat{\psi}_{\alpha}(t)\right| \frac{d H(t)}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle}{E_{\beta}(t)-E_{\alpha}(t)} . \tag{4.183}
\end{equation*}
$$

Without the implied $\lambda$ perturbation of $\left|\hat{\psi}_{\alpha}(t)\right\rangle$ we can from eq. (4.174) write the exact generalization of eq. (4.178) as

$$
\begin{equation*}
\frac{d}{d t} \bar{b}_{\alpha}(t)=-\sum_{\beta \neq \alpha} \bar{b}_{\beta}(t) e^{-\frac{i}{\hbar} \int_{0}^{t}\left(E_{\beta \alpha}\left(t^{\prime}\right)-\hbar \Gamma_{\beta \alpha}\left(t^{\prime}\right)\right) d t^{\prime}} \frac{\left\langle\hat{\psi}_{\alpha}(t)\right| \frac{d H(t)}{d t}\left|\hat{\psi}_{\beta}(t)\right\rangle}{E_{\beta}(t)-E_{\alpha}(t)} \tag{4.184}
\end{equation*}
$$

### 5.1 RECAP. WHERE WE GOT TO ON FERMI'S GOLDEN RULE.

We are continuing on the topic of Fermi golden rule, as also covered in $\S 17.2$ of the text [4]. Utilizing a wave train with peaks separation $\Delta t=2 \pi / \omega_{0}$, zero before some initial time fig. 5.1. Perturbing a state


Figure 5.1: Sine only after an initial time.
in the $i$ th energy level, and looking at the states for the $m$ th energy level as illustrated in fig. 5.2. Our matrix element was


Figure 5.2: Perturbation from $i$ to $m$ th energy levels.

$$
\begin{align*}
H_{m i}^{\prime}(t) & =2 A_{m} i \sin \left(\omega_{0} t\right) \theta(t) \\
& =i A_{m i}\left(e^{-i \omega_{0} t}-e^{i \omega_{0} t}\right) \theta(t) \tag{5.1}
\end{align*}
$$

and we found

$$
\begin{equation*}
c_{m}^{(1)}(t)=\frac{A_{m i}}{\hbar} \int_{0}^{t} d t^{\prime}\left(e^{i\left(\omega_{m i}-\omega_{0}\right) t^{\prime}}-e^{i\left(\omega_{m i}+\omega_{0}\right) t^{\prime}}\right) \tag{5.2}
\end{equation*}
$$

and argued that

$$
\begin{equation*}
\left|c_{m}^{(1)}(t)\right|^{2} \sim\left(\frac{A_{m i}}{\hbar}\right)^{2} t^{2}+\cdots \tag{5.3}
\end{equation*}
$$

where $\omega_{0} t \gg 1$ for $\omega_{m i} \sim \pm \omega_{0}$. We can also just integrate eq. (5.2) directly

$$
\begin{align*}
c_{m}^{(1)}(t) & =\frac{A_{m i}}{\hbar}\left(\frac{e^{i\left(\omega_{m i}-\omega_{0}\right) t}-1}{i\left(\omega_{m i}-\omega_{0}\right)}-\frac{e^{i\left(\omega_{m i}+\omega_{0}\right) t}-1}{i\left(\omega_{m i}+\omega_{0}\right)}\right)  \tag{5.4}\\
& \equiv A_{m i}\left(\omega_{0}, t\right)-A_{m i}\left(-\omega_{0}, t\right),
\end{align*}
$$

where

$$
\begin{equation*}
A_{m i}\left(\omega_{0}, t\right)=\frac{A_{m i}}{\hbar} \frac{e^{i\left(\omega_{m i}-\omega_{0}\right) t}-1}{i\left(\omega_{m i}-\omega_{0}\right)} . \tag{5.5}
\end{equation*}
$$

Factoring out the phase term, we have

$$
\begin{equation*}
A_{m i}\left(\omega_{0}, t\right)=\frac{A_{m i}}{\hbar} e^{i\left(\omega_{m i}-\omega_{0}\right) t / 2} \frac{\sin \left(\left(\omega_{m i}-\omega_{0}\right) t / 2\right)}{\left(\omega_{m i}-\omega_{0}\right) / 2} \tag{5.6}
\end{equation*}
$$

We will have two lobes, centered on $\pm \omega_{0}$, as illustrated in fig. 5.3.


Figure 5.3: Two sinc lobes.

## 5.2 fermi's golden rule.

Fermi's Golden rule applies to a continuum of states (there are other forms of Fermi's golden rule, but this is the one we will talk about, and is the one in the book). One example is the ionized states of an atom, where the energy level separation becomes so small that we can consider it continuous. Another example are the unbound states in a semiconductor well as illus-

Figure 5.4: Continuum of energy levels for ionized states of an atom.


Figure 5.5: Semi-conductor well.
trated in fig. 5.5. Note that we can have reflection from the well even in the continuum states where we would have no such reflection classically. However, with enough energy, states are approximately plane waves. In one dimension

$$
\begin{align*}
\left\langle x \mid \psi_{p}\right\rangle & \approx \frac{e^{i p x / \hbar}}{\sqrt{2 \pi \hbar}}  \tag{5.7}\\
\left\langle\psi_{p} \mid \psi_{p}^{\prime}\right\rangle & =\delta\left(p-p^{\prime}\right)
\end{align*}
$$

or in 3d

$$
\begin{align*}
\left\langle\mathbf{r} \mid \psi_{\mathbf{p}}\right\rangle & \approx \frac{e^{i \mathbf{p} \cdot \mathbf{r} / \hbar}}{(2 \pi \hbar)^{3 / 2}}  \tag{5.8}\\
\left\langle\psi_{\mathbf{p}} \mid \psi_{\mathbf{p}^{\prime}}\right\rangle & =\delta^{3}\left(\mathbf{p}-\mathbf{p}^{\prime}\right)
\end{align*}
$$

Let us consider the 1 d model for the quantum well in more detail. Including both discrete and continuous states we have

$$
\begin{equation*}
|\psi(t)\rangle=\sum_{n} c_{n}(t) e^{-i \omega_{n} t}\left|\psi_{n}\right\rangle+\int d p c_{p}(t) e^{-i \omega_{p} t}\left|\psi_{p}\right\rangle \tag{5.9}
\end{equation*}
$$

Imagine at $t=0$ that the wave function started in some discrete state, and look at the probability that we "kick the electron out of the well". Calculate

$$
\begin{equation*}
\mathcal{P}=\int d p\left|c_{p}^{(1)}(t)\right|^{2} \tag{5.10}
\end{equation*}
$$

Now, we assume that our matrix element has the following form

$$
\begin{equation*}
H_{p i}^{\prime}(t)=\left(\bar{A}_{p i} e^{-i \omega_{0} t}+\bar{B}_{p i} e^{i \omega_{0} t}\right) \theta(t) \tag{5.11}
\end{equation*}
$$

generalizing the wave train matrix element that we had previously

$$
\begin{equation*}
H_{m i}^{\prime}(t)=i A_{m i}\left(e^{-i \omega_{0} t}-e^{i \omega_{0} t}\right) \theta(t) \tag{5.12}
\end{equation*}
$$

Doing the perturbation we have

$$
\begin{equation*}
\mathcal{P}=\int d p\left|A_{p i}\left(\omega_{0}, t\right)+B_{p i}\left(-\omega_{0}, t\right)\right|^{2} \tag{5.13}
\end{equation*}
$$

where

$$
\begin{equation*}
A_{p i}\left(\omega_{0}, t\right)=\frac{2 \bar{A}_{p i}}{i \hbar} e^{i\left(\omega_{p i}-\omega_{0}\right) t / 2} \frac{\sin \left(\left(\omega_{p i}-\omega_{0}\right) t / 2\right)}{\omega_{p i}-\omega_{0}} \tag{5.14}
\end{equation*}
$$

which is peaked at $\omega_{p i}=\omega_{0}$, and

$$
\begin{equation*}
B_{p i}\left(\omega_{0}, t\right)=\frac{2 \bar{B}_{p i}}{i \hbar} e^{i\left(\omega_{p i}+\omega_{0}\right) t / 2} \frac{\sin \left(\left(\omega_{p i}+\omega_{0}\right) t / 2\right)}{\omega_{p i}+\omega_{0}} \tag{5.15}
\end{equation*}
$$

which is peaked at $\omega_{p i}=-\omega_{0}$.
In eq. (5.13) at $t \gg 0$ the only significant contribution is from the A portion as illustrated in fig. 5.6 where we are down in the wiggles of $A_{p i}$. Our probability to find the particle in the continuum range is now approximately

$$
\begin{equation*}
\mathcal{P}=\int d p\left|A_{p i}\left(\omega_{0}, t\right)\right|^{2} \tag{5.16}
\end{equation*}
$$

With

$$
\begin{equation*}
\omega_{p i}-\omega_{0}=\frac{1}{\hbar}\left(\frac{p^{2}}{2 m}-E_{i}\right)-\omega_{0} \tag{5.17}
\end{equation*}
$$

define $\bar{p}$ so that

$$
\begin{equation*}
0=\frac{1}{\hbar}\left(\frac{\bar{p}^{2}}{2 m}-E_{i}\right)-\omega_{0} \tag{5.18}
\end{equation*}
$$



Figure 5.6: .


Figure 5.7: Momentum space view.

In momentum space, we know have the sinc functions peaked at $\pm \bar{p}$ as in fig. 5.7. The probability that the electron goes to the right is then

$$
\begin{align*}
\mathcal{P}_{+} & =\int_{0}^{\infty} d p\left|c_{p}^{(1)}(t)\right|^{2} \\
& =\int_{0}^{\infty} d p\left|\bar{A}_{p i}\right|^{2} \frac{\sin ^{2}\left(\left(\omega_{p i}-\omega_{0}\right) t / 2\right)}{\left(\omega_{p i}-\omega_{0}\right)^{2}} \tag{5.19}
\end{align*}
$$

with

$$
\begin{equation*}
\omega_{p i}=\frac{1}{\hbar}\left(\frac{p^{2}}{2 m}-E_{i}\right), \tag{5.20}
\end{equation*}
$$

we have with a change of variables

$$
\begin{equation*}
\mathcal{P}_{+}=\frac{4}{\hbar^{2}} \int_{-E_{i} / \hbar}^{\infty} d \omega_{p i}\left|\bar{A}_{p i}\right|^{2} \frac{d p}{d \omega_{p i}} \frac{\sin ^{2}\left(\left(\omega_{p i}-\omega_{0}\right) t / 2\right)}{\left(\omega_{p i}-\omega_{0}\right)^{2}} \tag{5.21}
\end{equation*}
$$

Now suppose we have $t$ small enough so that $\mathcal{P}_{+} \ll 1$ and $t$ large enough SO

$$
\begin{equation*}
\left|\bar{A}_{p i}\right|^{2} \frac{d p}{d \omega_{p i}} \tag{5.22}
\end{equation*}
$$

is roughly constant over $\Delta \omega$. This is a sort of "Goldilocks condition", a time that can not be too small, and can not be too large, but instead has to be "just right". Given such a condition

$$
\begin{equation*}
\mathcal{P}_{+}=\frac{4}{\hbar^{2}}\left|\bar{A}_{p i}\right|^{2} \frac{d p}{d \omega_{p i}} \int_{-E_{i} / \hbar}^{\infty} d \omega_{p i} \frac{\sin ^{2}\left(\left(\omega_{p i}-\omega_{0}\right) t / 2\right)}{\left(\omega_{p i}-\omega_{0}\right)^{2}} \tag{5.23}
\end{equation*}
$$

where we can pull stuff out of the integral since the main contribution is at the peak. Provided $\bar{p}$ is large enough, using eq. (5.35), then

$$
\begin{align*}
\int_{-E_{i} / \hbar}^{\infty} d \omega_{p i} \frac{\sin ^{2}\left(\left(\omega_{p i}-\omega_{0}\right) t / 2\right)}{\left(\omega_{p i}-\omega_{0}\right)^{2}} & \approx \int_{-\infty}^{\infty} d \omega_{p i} \frac{\sin ^{2}\left(\left(\omega_{p i}-\omega_{0}\right) t / 2\right)}{\left(\omega_{p i}-\omega_{0}\right)^{2}}  \tag{5.24}\\
& =\frac{t}{2} \pi
\end{align*}
$$

leaving the probability of the electron with a going right continuum state as

$$
\begin{gather*}
\text { matrix element } \\
\mathcal{P}_{+}=\frac{4}{\hbar^{2}} \underbrace{\left|\bar{A}_{p i}\right|^{2}}_{\text {density of states }} \underbrace{}_{\left.\frac{d p}{d \omega_{p i}}\right|_{\bar{p}}} \frac{t}{2} \pi . \tag{5.25}
\end{gather*}
$$

The $d p / d \omega_{p i}$ is something like "how many continuous states are associated with a transition from a discrete frequency interval."

We can also get this formally from eq. (5.23) with

$$
\begin{equation*}
\frac{\sin ^{2}\left(\left(\omega_{p i}-\omega_{0}\right) t / 2\right)}{\left(\omega_{p i}-\omega_{0}\right)^{2}} \rightarrow \frac{t}{2} \pi \delta\left(\omega_{p i}-\omega_{0}\right) \tag{5.26}
\end{equation*}
$$

So

$$
\begin{align*}
c_{p}^{(1)}(t) & \rightarrow \frac{2 \pi t}{\hbar^{2}}\left|\bar{A}_{p i}\right|^{2} \delta\left(\omega_{p i}-\omega_{0}\right)  \tag{5.27}\\
& =\frac{2 \pi t}{\hbar}\left|\bar{A}_{p i}\right|^{2} \delta\left(E_{p i}-\hbar \omega_{0}\right)
\end{align*}
$$

where $\delta(a x)=\delta(x) /|a|$ has been used to pull in a factor of $\hbar$ into the delta.
The ratio of the coefficient to time is then

$$
\begin{equation*}
\frac{c_{p}^{(1)}(t)}{t}=\frac{2 \pi}{\hbar}\left|\bar{A}_{p i}\right|^{2} \delta\left(E_{p i}-\hbar \omega_{0}\right) . \tag{5.28}
\end{equation*}
$$

or "between friends"

$$
\begin{equation*}
\prime \frac{d c_{p}^{(1)}(t)^{\prime \prime}}{d t}=\frac{2 \pi}{\hbar}\left|\bar{A}_{p i}\right|^{2} \delta\left(E_{p i}-\hbar \omega_{0}\right) \tag{5.29}
\end{equation*}
$$

roughly speaking we have a "rate" or transitions from the discrete into the continuous. Here "rate" is in quotes since it does not hold for small t.
This has been worked out for $\mathcal{P}_{+}$. This can also be done for $\mathcal{P}_{-}$, the probability that the electron will end up in a left trending continuum state.

While the above is not a formal derivation, but illustrates the form of what is called Fermi's golden rule. Namely that such a rate has the structure

$$
\begin{equation*}
\frac{2 \pi}{\hbar} \times(\text { matrix element })^{2} \times \text { energy conservation. } \tag{5.30}
\end{equation*}
$$

## 5.3 problems.

Exercise 5.1 Sinc squared integral.
In the Fermi's golden rule lecture we used the result for the integral of the squared sinc function, that is

$$
\int_{-\infty}^{\infty} \frac{\sin ^{2}(x|\mu|)}{x^{2}} d x=\pi|\mu| .
$$

Prove this.
Answer for Exercise 5.1
We want to evaluate

$$
\begin{equation*}
\int_{-\infty}^{\infty} \frac{\sin ^{2}(x|\mu|)}{x^{2}} d x \tag{5.31}
\end{equation*}
$$

for which we must remind ourselves about the appropriate contours required to do the evaluation. We make a few change of variables

$$
\begin{align*}
\int_{-\infty}^{\infty} \frac{\sin ^{2}(x|\mu|)}{x^{2}} d x & =|\mu| \int_{-\infty}^{\infty} \frac{\sin ^{2}(y)}{y^{2}} d y \\
& =-i|\mu| \int_{-\infty}^{\infty} \frac{\left(e^{i y}-e^{-i y}\right)^{2}}{(2 i y)^{2}} i d y  \tag{5.32}\\
& =-\frac{i|\mu|}{4} \int_{-i \infty}^{i \infty} \frac{e^{2 z}+e^{-2 z}-2}{z^{2}} d z .
\end{align*}
$$



Figure 5.8: Contour distorted to one side of the double pole at the origin.

Now we pick a contour that is distorted to one side of the origin as in fig. 5.8. We employ Jordan's theorem ( $\$ 8.12$ [11]) now to pick the contours for each of the integrals since we need to ensure the $e^{ \pm z}$ terms converges as $R \rightarrow \infty$ for the $z=R e^{i \theta}$ part of the contour. We can write

$$
\begin{equation*}
\int_{-\infty}^{\infty} \frac{\sin ^{2}(x|\mu|)}{x^{2}} d x=-\frac{i|\mu|}{4}\left(\int_{C_{0}+C_{2}} \frac{e^{2 z}}{z^{2}} d z+\int_{C_{0}+C_{1}} \frac{e^{-2 z}}{z^{2}} d z-\int_{C_{0}+C_{1}} \frac{2}{z^{2}} d z\right) . \tag{5.33}
\end{equation*}
$$

The second two integrals both surround no poles (the contours were strategically placed to cowardly avoid the pole of the integrand at the origin), so we have only the first to deal with

$$
\begin{align*}
\int_{C_{0}+C_{2}} \frac{e^{2 z}}{z^{2}} d z & =\left.2 \pi i \frac{1}{1!} \frac{d}{d z} e^{2 z}\right|_{z=0}  \tag{5.34}\\
& =4 \pi i
\end{align*}
$$

Putting everything back together we have

$$
\begin{equation*}
\int_{-\infty}^{\infty} \frac{\sin ^{2}(x|\mu|)}{x^{2}} d x=-\frac{i|\mu|}{4} 4 \pi i=\pi|\mu| . \tag{5.35}
\end{equation*}
$$

On the cavalier choice of contours. The choice of which contours to pick above may seem pretty arbitrary, but they are for good reason. Sup-
pose you picked $C_{0}+C_{1}$ for the first integral. On the big $C_{1}$ arc, then with $\mathrm{a} z=R e^{i \theta}$ substitution we have

$$
\begin{align*}
\left|\int_{C_{1}} \frac{e^{2 z}}{z^{2}} d z\right| & =\left|\int_{\theta=\pi / 2}^{-\pi / 2} \frac{e^{2 R(\cos \theta+i \sin \theta)}}{R^{2} e^{2 i \theta}} R i e^{i \theta} d \theta\right| \\
& =\frac{1}{R}\left|\int_{\theta=\pi / 2}^{-\pi / 2} e^{2 R(\cos \theta+i \sin \theta)} e^{-i \theta} d \theta\right|  \tag{5.36}\\
& \leq \frac{1}{R} \int_{\theta=-\pi / 2}^{\pi / 2}\left|e^{2 R \cos \theta}\right| d \theta \\
& \leq \frac{\pi e^{2 R}}{R} .
\end{align*}
$$

This clearly doesn't have the zero convergence property that we desire. We need to pick the $C_{2}$ contour for the first (positive exponent) integral since in that $[\pi / 2,3 \pi / 2]$ range, $\cos \theta$ is always negative. We can however, use the $C_{1}$ contour for the second (negative exponent) integral. Explicitly, again by example, using $C_{2}$ contour for the first integral, over that portion of the arc we have

$$
\begin{align*}
\left|\int_{C_{2}} \frac{e^{2 z}}{z^{2}} d z\right| & =\left|\int_{\theta=\pi / 2}^{3 \pi / 2} \frac{e^{2 R(\cos \theta+i \sin \theta)}}{R^{2} e^{2 i \theta}} R i e^{i \theta} d \theta\right| \\
& =\frac{1}{R}\left|\int_{\theta=\pi / 2}^{3 \pi / 2} e^{2 R(\cos \theta+i \sin \theta)} e^{-i \theta} d \theta\right| \\
& \leq \frac{1}{R} \int_{\theta=\pi / 2}^{3 \pi / 2}\left|e^{2 R \cos \theta} d \theta\right|  \tag{5.37}\\
& \approx \frac{1}{R} \int_{\theta=\pi / 2}^{3 \pi / 2}\left|e^{-2 R} d \theta\right| \\
& =\frac{\pi e^{-2 R}}{R}
\end{align*}
$$

WKB METHOD.

## 6.1 wкb (wentzel-kramers-brillouin) method.

This lecture covers material from $\S 24$ of the text [4], and $\S 8$ of [8].
We start with the 1D time independent Schrödinger equation

$$
\begin{equation*}
-\frac{\hbar^{2}}{2 m} \frac{d^{2} U}{d x^{2}}+V(x) U(x)=E U(x) \tag{6.1}
\end{equation*}
$$

which we can write as

$$
\begin{equation*}
\frac{d^{2} U}{d x^{2}}+\frac{2 m}{\hbar^{2}}(E-V(x)) U(x)=0 . \tag{6.2}
\end{equation*}
$$

Consider a finite well potential as in fig. 6.1. With


Figure 6.1: Finite well potential.

$$
\begin{array}{ll}
k^{2}=\frac{2 m(E-V)}{\hbar}, & E>V \\
\kappa^{2}=\frac{2 m(V-E)}{\hbar}, & V>E, \tag{6.3}
\end{array}
$$

the bound state within the well is approximately

$$
\begin{equation*}
U \propto e^{ \pm i k x} \tag{6.4}
\end{equation*}
$$

and outside of the well is of the form

$$
\begin{equation*}
U \propto e^{ \pm \kappa x} . \tag{6.5}
\end{equation*}
$$

In general we can hope for something similar. Let us look for that something, but allow the constants $k$ and $\kappa$ to be functions of position

$$
\begin{array}{ll}
k^{2}(x)=\frac{2 m(E-V(x))}{\hbar}, & E>V \\
\kappa^{2}(x)=\frac{2 m(V(x)-E)}{\hbar}, & V>E \tag{6.6}
\end{array}
$$

In terms of $k$ Schrödinger's equation is just

$$
\begin{equation*}
\frac{d^{2} U(x)}{d x^{2}}+k^{2}(x) U(x)=0 \tag{6.7}
\end{equation*}
$$

We use the trial solution

$$
\begin{equation*}
U(x)=A e^{i \phi(x)} \tag{6.8}
\end{equation*}
$$

allowing $\phi(x)$ to be complex

$$
\begin{equation*}
\phi(x)=\phi_{R}(x)+i \phi_{I}(x) . \tag{6.9}
\end{equation*}
$$

We need second derivatives

$$
\begin{align*}
\left(e^{i \phi}\right)^{\prime \prime} & =\left(i \phi^{\prime} e^{i \phi}\right)^{\prime} \\
& =\left(i \phi^{\prime}\right)^{2} e^{i \phi}+i \phi^{\prime \prime} e^{i \phi}, \tag{6.10}
\end{align*}
$$

and plug back into our Schrödinger equation to obtain

$$
\begin{equation*}
-\left(\phi^{\prime}(x)\right)^{2}+i \phi^{\prime \prime}(x)+k^{2}(x)=0 \tag{6.11}
\end{equation*}
$$

For the first round of approximation we assume

$$
\begin{equation*}
\phi^{\prime \prime}(x) \approx 0, \tag{6.12}
\end{equation*}
$$

and obtain

$$
\begin{equation*}
\left(\phi^{\prime}(x)\right)^{2}=k^{2}(x), \tag{6.13}
\end{equation*}
$$

or

$$
\begin{equation*}
\phi^{\prime}(x)= \pm k(x) . \tag{6.14}
\end{equation*}
$$

A second round of approximation we use eq. (6.14) and obtain

$$
\begin{equation*}
\phi^{\prime \prime}(x)= \pm k^{\prime}(x) . \tag{6.15}
\end{equation*}
$$

Plugging back into eq. (6.11) we have

$$
\begin{equation*}
-\left(\phi^{\prime}(x)\right)^{2} \pm i k^{\prime}(x)+k^{2}(x)=0 . \tag{6.16}
\end{equation*}
$$

Things get a little confusing here with the $\pm$ variation since we have to take a second set of square roots, so let's consider these separately.

Case I. positive root. With $\phi^{\prime}=+k$, we have

$$
\begin{equation*}
-\left(\phi^{\prime}(x)\right)^{2}+i k^{\prime}(x)+k^{2}(x)=0 \tag{6.17}
\end{equation*}
$$

or

$$
\begin{align*}
\phi^{\prime}(x) & = \pm \sqrt{+i k^{\prime}(x)+k^{2}(x)} \\
& = \pm k(x) \sqrt{1+i \frac{k^{\prime}(x)}{k^{2}(x)}} \tag{6.18}
\end{align*}
$$

If $k^{\prime}$ is small compared to $k^{2}$

$$
\begin{equation*}
\frac{k^{\prime}(x)}{k^{2}(x)} \ll 1 \tag{6.19}
\end{equation*}
$$

then we have

$$
\begin{equation*}
\phi^{\prime}(x) \approx \pm k(x)\left(1+i \frac{k^{\prime}(x)}{2 k^{2}(x)}\right)= \pm\left(k(x)+i \frac{k^{\prime}(x)}{2 k(x)}\right) \tag{6.20}
\end{equation*}
$$

Since we'd picked $\phi^{\prime} \approx+k$ in this case, we pick the positive sign, and can now integrate

$$
\begin{align*}
\phi(x) & =\int d x k(x)+i \int d x \frac{k^{\prime}(x)}{2 k(x)}+\ln \text { const }  \tag{6.21}\\
& =\int d x k(x)+i \frac{1}{2} \ln k(x)+\ln \text { const. }
\end{align*}
$$

Going back to our wavefunction, for this $E>V(x)$ case we have

$$
\begin{align*}
U(x) & \sim e^{i \phi(x)} \\
& =\exp \left(i\left(\int d x k(x)+i \frac{1}{2} \ln k(x)+\mathrm{const}\right)\right)  \tag{6.22}\\
& \sim \exp \left(i\left(\int d x k(x)+i \frac{1}{2} \ln k(x)\right)\right) \\
& =e^{i \int d x k(x)} e^{-\frac{1}{2} \ln k(x)},
\end{align*}
$$

or

$$
\begin{equation*}
U(x) \propto \frac{1}{\sqrt{k(x)}} e^{i \int d x k(x)} \tag{6.23}
\end{equation*}
$$

Case II: negative sign. Now treat $\phi^{\prime} \approx-k$. This gives us

$$
\begin{equation*}
\phi^{\prime}(x) \approx \pm k(x)\left(1-i \frac{k^{\prime}(x)}{2 k^{2}(x)}\right)= \pm\left(k(x)-i \frac{k^{\prime}(x)}{2 k(x)}\right) . \tag{6.24}
\end{equation*}
$$

This time we want the negative root to match $\phi^{\prime} \approx-k$. Integrating, we have

$$
\begin{align*}
i \phi(x) & =-i \int d x\left(k(x)-i \frac{k^{\prime}(x)}{2 k(x)}\right) \\
& =-i \int k(x) d x-\frac{1}{2} \int \frac{k^{\prime}}{k} d x  \tag{6.25}\\
& =-i \int k(x) d x-\frac{1}{2} \ln k+\ln \text { constant. }
\end{align*}
$$

This gives us

$$
\begin{equation*}
U(x) \propto \frac{1}{\sqrt{k(x)}} e^{-i \int d x k(x)} \tag{6.26}
\end{equation*}
$$

Provided we have eq. (6.19), we can summarize these as

$$
\begin{equation*}
U(x) \propto \frac{1}{\sqrt{k(x)}} e^{ \pm i \int d x k(x)} \tag{6.27}
\end{equation*}
$$

It's not hard to show that for the $E<V(x)$ case we find

$$
\begin{equation*}
U(x) \propto \frac{1}{\sqrt{\kappa(x)}} e^{ \pm \int d x \kappa(x)} \tag{6.28}
\end{equation*}
$$

this time, provided that our potential satisfies

$$
\begin{equation*}
\frac{\kappa^{\prime}(x)}{\kappa^{2}(x)} \ll 1 . \tag{6.29}
\end{equation*}
$$

Validity. For this to be valid we require that $V(x)$ change very slowly. This implies that $k^{\prime}(x)$ is small, and $k(x)=\sqrt{2 m(E-V(x))} / \hbar$. We need $E$ very far away from the potential. In particular, $|(E-V(x)) / V(x)| \gg 1$.

## 6.2 turning points..

WKB will not work at the turning points in this figure since our main assumption was that

$$
\begin{equation*}
\left|\frac{k^{\prime}(x)}{k^{2}(x)}\right| \ll 1, \tag{6.30}
\end{equation*}
$$



Figure 6.2: Example of a general potential.


Figure 6.3: Turning points where WKB will not work.


Figure 6.4: Diagram for patching method discussion.
so we get into trouble where $k(x) \sim 0$. There are some methods for dealing with this. Our text as well as Griffiths give some examples, but they require Bessel functions and more complex mathematics.

The idea is that one finds the WKB solution in the regions of validity, and then looks for a polynomial solution in the patching region where we are closer to the turning point, probably requiring lookup of various special functions.

This power series method is also outlined in [22], where solutions to connect the regions are expressed in terms of Airy functions.

## 6.3 examples.

## Example 6.1: Infinite well potential.

Consider the potential

$$
V(x)= \begin{cases}v(x) & \text { if } x \in[0, a]  \tag{6.31}\\ \infty & \text { otherwise }\end{cases}
$$

as illustrated in fig. 6.5.


Figure 6.5: Arbitrary potential in an infinite well.
Inside the well, we have

$$
\begin{equation*}
\psi(x)=\frac{1}{\sqrt{k(x)}}\left(C_{+} e^{i \int_{0}^{x} k\left(x^{\prime}\right) d x^{\prime}}+C_{-} e^{-i \int_{0}^{x} k\left(x^{\prime}\right) d x^{\prime}}\right) \tag{6.32}
\end{equation*}
$$

where

$$
\begin{equation*}
k(x)=\frac{1}{\hbar} \sqrt{2 m(E-v(x)} \tag{6.33}
\end{equation*}
$$

## With

$$
\begin{equation*}
\phi(x)=e^{\int_{0}^{x} k\left(x^{\prime}\right) d x^{\prime}} \tag{6.34}
\end{equation*}
$$

We have

$$
\begin{align*}
\psi(x) & =\frac{1}{\sqrt{k(x)}}\left(C_{+}(\cos \phi+i \sin \phi)+C_{-}(\cos \phi-i \sin \phi)\right) \\
& =\frac{1}{\sqrt{k(x)}}\left(\left(C_{+}+C_{-}\right) \cos \phi+i\left(C_{+}-C_{-}\right) \sin \phi\right)  \tag{6.35}\\
& =\frac{1}{\sqrt{k(x)}}\left(\left(C_{+}+C_{-}\right) \cos \phi+i\left(C_{+}-C_{-}\right) \sin \phi\right) \\
& \equiv \frac{1}{\sqrt{k(x)}}\left(C_{2} \cos \phi+C_{1} \sin \phi\right),
\end{align*}
$$

Where

$$
\begin{align*}
& C_{2}=C_{+}+C_{-}  \tag{6.36}\\
& C_{1}=i\left(C_{+}-C_{-}\right)
\end{align*}
$$

Setting boundary conditions we have

$$
\begin{equation*}
\phi(0)=0 \tag{6.37}
\end{equation*}
$$

Noting that we have $\phi(0)=0$, we have

$$
\begin{equation*}
\frac{1}{\sqrt{k(0)}} C_{2}=0 \tag{6.38}
\end{equation*}
$$

So

$$
\begin{equation*}
\psi(x) \sim \frac{1}{\sqrt{k(x)}} \sin \phi \tag{6.39}
\end{equation*}
$$

At the other boundary

$$
\begin{equation*}
\psi(a)=0 \tag{6.40}
\end{equation*}
$$

So we require

$$
\begin{equation*}
\sin \phi(a)=\sin (n \pi) \tag{6.41}
\end{equation*}
$$

or

$$
\begin{equation*}
\frac{1}{\hbar} \int_{0}^{a} \sqrt{2 m\left(E-v\left(x^{\prime}\right)\right.} d x^{\prime}=n \pi \tag{6.42}
\end{equation*}
$$

This is called the Bohr-Sommerfeld condition.

Check, with $v(x)=0$. We have

$$
\begin{equation*}
\frac{1}{\hbar} \sqrt{2 m E} a=n \pi \tag{6.43}
\end{equation*}
$$

or

$$
\begin{equation*}
E=\frac{1}{2 m}\left(\frac{n \pi \hbar}{a}\right)^{2} \tag{6.44}
\end{equation*}
$$

## 7.1 hilbert spaces.

READING: $\S 30$ of the text [4] covers entangled states. The rest of the composite state background is buried somewhere in some of the advanced material sections. FIXME: what section?

Example, one spin one half particle and one spin one particle. We can describe either quantum mechanically, described by a pair of Hilbert spaces $H_{1}$, of dimension $D_{1}$, and $H_{2}$, of dimension $D_{2}$.

We loosely describe a Hilbert space (finite or infinite dimensional) as the set of states that describe the system. There were some additional details (completeness, normalizable, $L 2$ integrable, ...) not really covered in the physics curriculum, but available in mathematical descriptions.

We form the composite (Hilbert) space

$$
\begin{equation*}
H=H_{1} \otimes H_{2}, \tag{7.1}
\end{equation*}
$$

$$
\begin{equation*}
H_{1}:\left|\phi_{1}^{(i)}\right\rangle \tag{7.2}
\end{equation*}
$$

for any ket in $H_{1}$

$$
\begin{equation*}
|I\rangle=\sum_{i=1}^{D_{1}} c_{i}\left|\phi_{1}^{(i)}\right\rangle, \tag{7.3}
\end{equation*}
$$

where

$$
\begin{equation*}
\left\langle\phi_{1}^{(i)} \mid \phi_{1}^{(j)}\right\rangle=\delta^{i j} \tag{7.4}
\end{equation*}
$$

Similarly

$$
\begin{equation*}
H_{2}:\left|\phi_{2}^{(i)}\right\rangle \tag{7.5}
\end{equation*}
$$

for any ket in $\mathrm{H}_{2}$

$$
\begin{equation*}
|I I\rangle=\sum_{i=1}^{D_{2}} d_{i}\left|\phi_{2}^{(i)}\right\rangle \tag{7.6}
\end{equation*}
$$

where

$$
\begin{equation*}
\left\langle\phi_{2}^{(i)} \mid \phi_{2}^{(j)}\right\rangle=\delta^{i j} . \tag{7.7}
\end{equation*}
$$

The composite Hilbert space has dimension $D_{1} D_{2}$.
The basis kets are

$$
\begin{equation*}
\left|\phi_{1}^{(i)}\right\rangle \otimes\left|\phi_{2}^{(j)}\right\rangle=\left|\phi^{(i j)}\right\rangle, \tag{7.8}
\end{equation*}
$$

where

$$
\begin{equation*}
\left\langle\phi^{(i j)} \mid \phi^{(k l)}\right\rangle=\delta^{i k} \delta^{j l} . \tag{7.9}
\end{equation*}
$$

Any ket in $H$ can be written

$$
\begin{align*}
|\psi\rangle & =\sum_{i=1}^{D_{1}} \sum_{j=1}^{D_{2}} f_{i j}\left|\phi_{1}^{(i)}\right\rangle \otimes\left|\phi_{2}^{(j)}\right\rangle  \tag{7.10}\\
& =\sum_{i=1}^{D_{1}} \sum_{j=1}^{D_{2}} f_{i j}\left|\phi^{(i j)}\right\rangle .
\end{align*}
$$

Direct product of kets:

$$
\begin{align*}
|I\rangle \otimes|I I\rangle & \equiv \sum_{i=1}^{D_{1}} \sum_{j=1}^{D_{2}} c_{i} d_{j}\left|\phi_{1}^{(i)}\right\rangle \otimes\left|\phi_{2}^{(j)}\right\rangle \\
& =\sum_{i=1}^{D_{1}} \sum_{j=1}^{D_{2}} c_{i} d_{j}\left|\phi^{(i j)}\right\rangle . \tag{7.11}
\end{align*}
$$

If $|\psi\rangle$ in $H$ cannot be written as $|I\rangle \otimes|I I\rangle$, then $|\psi\rangle$ is said to be "entangled".
FIXME: insert a concrete example of this, with some low dimension.

## 7.2 operators.

With operators $O_{1}$ and $O_{2}$ on the respective Hilbert spaces. We would now like to build

$$
\begin{equation*}
O_{1} \otimes O_{2} \tag{7.12}
\end{equation*}
$$

If one defines

$$
\begin{equation*}
O_{1} \otimes O_{2} \equiv \sum_{i=1}^{D_{1}} \sum_{j=1}^{D_{2}} f_{i j}\left|O_{1} \phi_{1}^{(i)}\right\rangle \otimes\left|O_{2} \phi_{2}^{(j)}\right\rangle . \tag{7.13}
\end{equation*}
$$

Q:Can every operator that can be defined on the composite space have a representation of this form? No.

Special cases. The identity operators. Suppose that

$$
\begin{equation*}
|\psi\rangle=\sum_{i=1}^{D_{1}} \sum_{j=1}^{D_{2}} f_{i j}\left|\phi_{1}^{(i)}\right\rangle \otimes\left|\phi_{2}^{(j)}\right\rangle \tag{7.14}
\end{equation*}
$$

then

$$
\begin{equation*}
\left(O_{1} \otimes I_{2}\right)|\psi\rangle=\sum_{i=1}^{D_{1}} \sum_{j=1}^{D_{2}} f_{i j}\left|O_{1} \phi_{1}^{(i)}\right\rangle \otimes\left|\phi_{2}^{(j)}\right\rangle \tag{7.15}
\end{equation*}
$$

## Example 7.1: A commutator.

Can do other operations. Example:

$$
\begin{equation*}
\left[O_{1} \otimes I_{2}, I_{1} \otimes O_{2}\right]=0 \tag{7.16}
\end{equation*}
$$

Let us verify this one. Suppose that our state has the representation

$$
\begin{equation*}
|\psi\rangle=\sum_{i=1}^{D_{1}} \sum_{j=1}^{D_{2}} f_{i j}\left|\phi_{1}^{(i)}\right\rangle \otimes\left|\phi_{2}^{(j)}\right\rangle \tag{7.17}
\end{equation*}
$$

so that the action on this ket from the composite operations are

$$
\begin{align*}
& \left(O_{1} \otimes I_{2}\right)|\psi\rangle=\sum_{i=1}^{D_{1}} \sum_{j=1}^{D_{2}} f_{i j}\left|O_{1} \phi_{1}^{(i)}\right\rangle \otimes\left|\phi_{2}^{(j)}\right\rangle \\
& \left(I_{1} \otimes O_{2}\right)|\psi\rangle=\sum_{i=1}^{D_{1}} \sum_{j=1}^{D_{2}} f_{i j}\left|\phi_{1}^{(i)}\right\rangle \otimes\left|O_{2} \phi_{2}^{(j)}\right\rangle . \tag{7.18}
\end{align*}
$$

## Our commutator is

$$
\begin{align*}
& {\left[\left(O_{1} \otimes I_{2}\right),\left(I_{1} \otimes O_{2}\right)\right]|\psi\rangle } \\
= & \left(O_{1} \otimes I_{2}\right)\left(I_{1} \otimes O_{2}\right)|\psi\rangle-\left(I_{1} \otimes O_{2}\right)\left(O_{1} \otimes I_{2}\right)|\psi\rangle \\
= & \left(O_{1} \otimes I_{2}\right) \sum_{i=1}^{D_{1}} \sum_{j=1}^{D_{2}} f_{i j}\left|\phi_{1}^{(i)}\right\rangle \otimes\left|O_{2} \phi_{2}^{(j)}\right\rangle \\
& -\left(I_{1} \otimes O_{2}\right) \sum_{i=1}^{D_{1}} \sum_{j=1}^{D_{2}} f_{i j}\left|O_{1} \phi_{1}^{(i)}\right\rangle \otimes\left|\phi_{2}^{(j)}\right\rangle \\
= & \sum_{i=1}^{D_{1}} \sum_{j=1}^{D_{2}} f_{i j}\left|O_{1} \phi_{1}^{(i)}\right\rangle \otimes\left|O_{2} \phi_{2}^{(j)}\right\rangle-\sum_{i=1}^{D_{1}} \sum_{j=1}^{D_{2}} f_{i j}\left|O_{1} \phi_{1}^{(i)}\right\rangle \otimes\left|O_{2} \phi_{2}^{(j)}\right\rangle \\
= & 0 . \tag{7.19}
\end{align*}
$$

### 7.3 GENERALIZATIONS.

Can generalize to

$$
\begin{equation*}
H_{1} \otimes H_{2} \otimes H_{3} \otimes \cdots \tag{7.20}
\end{equation*}
$$

Can also start with $H$ and seek factor spaces. If $H$ is not prime there are, in general, many ways to find factor spaces

$$
\begin{equation*}
H=H_{1} \otimes H_{2}=H_{1}^{\prime} \otimes H_{2}^{\prime} . \tag{7.21}
\end{equation*}
$$

A ket $|\psi\rangle$, if unentangled in the first factor space, then it will be in general entangled in a second space. Thus ket entanglement is not a property of the ket itself, but instead is intrinsically related to the space in which it is represented.

## 7.4 recalling the stern-gerlach system from phy 354 .

We had one example of a composite system in phy356 that I recall. It was related to states of the silver atoms in a Stern Gerlach apparatus, where we had one state from the Hamiltonian that governs position and momentum and another from the Hamiltonian for the spin, where each of these states was considered separately.

This makes me wonder what would the Hamiltonian for a system (say a single electron) that includes both spin and position/momentum would look like, and how is it that one can solve this taking spin and non-spin states separately?

Professor Sipe, when asked said of this
"It is complicated because not only would the spin of the electron interact with the magnetic field, but its translational motion would respond to the magnetic field too. A simpler case is a neutral atom with an electron with an unpaired spin. Then there is no Lorentz force on the atom itself. The Hamiltonian is just the sum of a free particle Hamiltonian and a Zeeman term due to the spin interacting with the magnetic field. This is precisely the Stern-Gerlach problem."

I did not remember what the Zeeman term looked like, but wikipedia does [23], and it is the magnetic field interaction

$$
\begin{equation*}
-\boldsymbol{\mu} \cdot \mathbf{B}, \tag{7.22}
\end{equation*}
$$

that we get when we gauge transform the Dirac equation for the electron as covered in $\S 36.4$ of the text (also introduced in chapter 6, which was not covered in class). That does not look too much like how we studied the Stern-Gerlach problem? I thought that for that problem we had a Hamiltonian of the form

$$
\begin{equation*}
H=a_{i j}|i\rangle\langle j| . \tag{7.23}
\end{equation*}
$$

It is not clear to me how this ket-bra Hamiltonian and the Zeeman Hamiltonian are related (ie: the spin Hamiltonians that we used in 356 and were on old 356 exams were all pulled out of magic hats and it was not obvious where these came from).

### 8.1 Generators.

Covered in $\S 26$ of the text [4].

## Example 8.1: Time translation.

One of the simplest examples of a generator, is that of time translation

$$
\begin{equation*}
|\psi(t)\rangle=e^{-i H t / \hbar}|\psi(0)\rangle . \tag{8.1}
\end{equation*}
$$

The Hamiltonian "generates" evolution (or translation) in time.

## Example 8.2: Spatial translation.

Similarily, we may apply an operator to translate the ket associated with a given position

$$
\begin{equation*}
|\mathbf{r}+\mathbf{a}\rangle=e^{-i \mathbf{a} \cdot \mathbf{P} / \hbar}|\mathbf{r}\rangle \tag{8.2}
\end{equation*}
$$



Figure 8.1: Vector translation.
$\mathbf{P}$ is the operator that generates translations. Written out, we have

$$
\begin{align*}
e^{-i \mathbf{a} \cdot \mathbf{P} / \hbar} & =e^{-i\left(a_{x} P_{x}+a_{y} P_{y}+a_{z} P_{z}\right) / \hbar} \\
& =e^{-i a_{x} P_{x} / \hbar} e^{-i a_{y} P_{y} / \hbar} e^{-i a_{z} P_{z} / \hbar} \tag{8.3}
\end{align*}
$$

where the factorization was possible because $P_{x}, P_{y}$, and $P_{z}$ commute

$$
\begin{equation*}
\left[P_{i}, P_{j}\right]=0 \tag{8.4}
\end{equation*}
$$

for any $i, j$ (including $i=i$ as I dumbly questioned in class $\ldots$ this is a commutator, so $\left[P_{i}, P_{j}\right]=P_{i} P_{i}-P_{i} P_{i}=0$ ).

The fact that the $P_{i}$ commute means that successive translations can be done in any order and have the same result.

In class we were rewarded with a graphic demo of translation component commutation as Professor Sipe pulled a giant wood carving of a cat (or tiger?) out from beside the desk and proceeded to translate it around on the desk in two different orders, with the cat ending up in the same place each time.

Exponential commutation. Note that in general

$$
\begin{equation*}
e^{A+B} \neq e^{A} e^{B} \tag{8.5}
\end{equation*}
$$

unless $[A, B]=0$. To show this one can compare

$$
\begin{align*}
e^{A+B} & =1+A+B+\frac{1}{2}(A+B)^{2}+\cdots  \tag{8.6}\\
& =1+A+B+\frac{1}{2}\left(A^{2}+A B+B A+B^{2}\right)+\cdots
\end{align*}
$$

and

$$
\begin{align*}
e^{A} e^{B} & =\left(1+A+\frac{1}{2} A^{2}+\cdots\right)\left(1+B+\frac{1}{2} B^{2}+\cdots\right)  \tag{8.7}\\
& =1+A+B+\frac{1}{2}\left(A^{2}+2 A B+B^{2}\right)+\cdots
\end{align*}
$$

Comparing the second order (for example) we see that we must have for equality

$$
\begin{equation*}
A B+B A=2 A B \tag{8.8}
\end{equation*}
$$

or

$$
\begin{equation*}
B A=A B \tag{8.9}
\end{equation*}
$$

or

$$
\begin{equation*}
[A, B]=0 \tag{8.10}
\end{equation*}
$$

Translating a ket. If we consider the quantity

$$
\begin{equation*}
e^{-i a \cdot \mathbf{P} / \hbar}|\psi\rangle=\left|\psi^{\prime}\right\rangle, \tag{8.11}
\end{equation*}
$$

does this ket "translated" by a make any sense? The vector a lives in a 3D space and our ket $|\psi\rangle$ lives in Hilbert space. A quantity like this deserves some careful thought and is the subject of some such thought in the Interpretations of Quantum mechanics course. For now, we can think of the operator and ket as a "gadget" that prepares a state.

A student in class pointed out that $|\psi\rangle$ can be dependent on many degrees of freedom, for example, the positions of eight different particles. This translation gadget in such a case acts on the whole kit and caboodle.

Now consider the matrix element

$$
\begin{equation*}
\left\langle\mathbf{r} \mid \psi^{\prime}\right\rangle=\langle\mathbf{r}| e^{-i \mathbf{a} \cdot \mathbf{P} / \hbar}|\psi\rangle . \tag{8.12}
\end{equation*}
$$

Note that

$$
\begin{align*}
\langle\mathbf{r}| e^{-i \mathbf{a} \cdot \mathbf{P} / \hbar} & =\left(e^{i \mathbf{a} \cdot \mathbf{P} / \hbar}|\mathbf{r}\rangle\right)^{\dagger}  \tag{8.13}\\
& =(|\mathbf{r}-\mathbf{a}\rangle)^{\dagger},
\end{align*}
$$

so

$$
\begin{equation*}
\left\langle\mathbf{r} \mid \psi^{\prime}\right\rangle=\langle\mathbf{r}-\mathbf{a} \mid \psi\rangle, \tag{8.14}
\end{equation*}
$$

or

$$
\begin{equation*}
\psi^{\prime}(\mathbf{r})=\psi(\mathbf{r}-\mathbf{a}) . \tag{8.15}
\end{equation*}
$$

This is what we expect of a translated function, as illustrated in fig. 8.2.


Figure 8.2: Active spatial translation.

## Example 8.3: Spatial rotation.

We have been introduced to the angular momentum operator

$$
\begin{equation*}
\mathbf{L}=\mathbf{R} \times \mathbf{P} \tag{8.16}
\end{equation*}
$$

where

$$
\begin{align*}
L_{x} & =Y P_{z}-Z P_{y} \\
L_{y} & =Z P_{x}-X P_{z}  \tag{8.17}\\
L_{z} & =X P_{y}-Y P_{x}
\end{align*}
$$

We also found that

$$
\begin{equation*}
\left[L_{i}, L_{j}\right]=i \hbar \sum_{k} \epsilon_{i j k} L_{k} \tag{8.18}
\end{equation*}
$$

These non-zero commutators show that the components of angular momentum do not commute.

Define

$$
\begin{equation*}
|\mathcal{R}(\mathbf{r})\rangle=e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{L} / \hbar}|\mathbf{r}\rangle \tag{8.19}
\end{equation*}
$$

This is the vector that we get by actively rotating the vector $\mathbf{r}$ by an angle $\theta$ counterclockwise about $\hat{\mathbf{n}}$, as in fig. 8.3.


Figure 8.3: Active vector rotations.
An active rotation rotates the vector, leaving the coordinate system fixed, whereas a passive rotation is one for which the coordinate system is rotated, and the vector is left fixed.

Note that rotations do not commute. Suppose that we have a pair of rotations as in fig. 8.4.


Figure 8.4: A example pair of non-commuting rotations.
Again, we get the graphic demo, with Professor Sipe rotating the big wooden cat sculpture. Did he bring that in to class just to make this point (too bad I missed the first couple minutes of the lecture).

Rather amusingly, he points out that most things in life do not commute. We get much different results if we apply the operations of putting water into the teapot and turning on the stove in different orders.

Rotating a ket. With a rotation gadget

$$
\begin{equation*}
\left|\psi^{\prime}\right\rangle=e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{L} / \hbar}|\psi\rangle, \tag{8.20}
\end{equation*}
$$

we can form the matrix element

$$
\begin{equation*}
\left\langle\mathbf{r} \mid \psi^{\prime}\right\rangle=\langle\mathbf{r}| e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{L} / \hbar}|\psi\rangle \tag{8.21}
\end{equation*}
$$

In this we have

$$
\begin{align*}
\langle\mathbf{r}| e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{L} / \hbar} & =\left(e^{i \theta \hat{\mathbf{n}} \cdot \mathbf{L} / \hbar}|\mathbf{r}\rangle\right)^{\dagger}  \tag{8.22}\\
& =\left(\left|\mathcal{R}^{-1}(\mathbf{r})\right\rangle\right)^{\dagger}
\end{align*}
$$

So

$$
\begin{equation*}
\left\langle\mathbf{r} \mid \psi^{\prime}\right\rangle=\left\langle\mathcal{R}^{-1}(\mathbf{r}) \mid \psi^{\prime}\right\rangle \tag{8.23}
\end{equation*}
$$

or

$$
\begin{equation*}
\psi^{\prime}(\mathbf{r})=\psi\left(\mathcal{R}^{-1}(\mathbf{r})\right) \tag{8.24}
\end{equation*}
$$

### 8.2 GENERALIZATIONS.

Recall what you did last year, where $H, \mathbf{P}$, and $\mathbf{L}$ were defined mechanically. We found

- $H$ generates time evolution (or translation in time).
- P generates spatial translation.
- L generates spatial rotation.

For our mechanical definitions we have

$$
\begin{equation*}
\left[P_{i}, P_{j}\right]=0 \tag{8.25}
\end{equation*}
$$

and

$$
\begin{equation*}
\left[L_{i}, L_{j}\right]=i \hbar \sum_{k} \epsilon_{i j k} L_{k} \tag{8.26}
\end{equation*}
$$

These are the relations that show us the way translations and rotations combine. We want to move up to a higher plane, a new level of abstraction. To do so we define $H$ as the operator that generates time evolution. If we have a theory that covers the behavior of how anything evolves in time, $H$ encodes the rules for this time evolution.
$\underline{\text { Define } \mathbf{P}}$ as the operator that generates translations in space.

Define $\mathbf{J}$ as the operator that generates rotations in space.
In order that these match expectations, we require

$$
\begin{equation*}
\left[P_{i}, P_{j}\right]=0, \tag{8.27}
\end{equation*}
$$

and

$$
\begin{equation*}
\left[J_{i}, J_{j}\right]=i \hbar \sum_{k} \epsilon_{i j k} J_{k} . \tag{8.28}
\end{equation*}
$$

In the simple theory of a spin less particle we have

$$
\begin{equation*}
\mathbf{J} \equiv \mathbf{L}=\mathbf{R} \times \mathbf{P} \tag{8.29}
\end{equation*}
$$

We actually need a generalization of this since this is, in fact, not good enough, even for low energy physics.

Many component wave functions. We are free to construct tuples of spatial vector functions like

$$
\left[\begin{array}{c}
\Psi_{I}(\mathbf{r}, t)  \tag{8.30}\\
\Psi_{I I}(\mathbf{r}, t)
\end{array}\right],
$$

or

$$
\left[\begin{array}{c}
\Psi_{I}(\mathbf{r}, t)  \tag{8.31}\\
\Psi_{I I}(\mathbf{r}, t) \\
\Psi_{I I I}(\mathbf{r}, t)
\end{array}\right],
$$

etc.
We will see that these behave qualitatively different than one component wave functions. We also do not have to be considering multiple particle wave functions, but just one particle that requires three functions in $\mathbb{R}^{3}$ to describe it (ie: we are moving in on spin).

Question: Do these live in the same vector space?
Answer: We will get to this.

A classical analogy. "There is only bad analogies, since if the are good they would be describing the same thing. We can however, produce some useful bad analogies"

1. A temperature field $T(\mathbf{r})$.
2. Electric field $\left(E_{x}(\mathbf{r}), E_{y}(\mathbf{r}), E_{z}(\mathbf{r})\right.$.

These behave in a much different way. If we rotate a scalar field like $T(\mathbf{r})$ as in fig. 8.5. Suppose we have a temperature field generated by, say,


Figure 8.5: Rotated temperature (scalar) field.
a match. Rotating the match above, we have

$$
\begin{equation*}
T^{\prime}(\mathbf{r})=T\left(\mathcal{R}^{-1}(\mathbf{r})\right) . \tag{8.32}
\end{equation*}
$$

Compare this to the rotation of an electric field, perhaps one produced by a capacitor, as in fig. 8.6. Is it true that we have


Figure 8.6: Rotating a capacitance electric field.

$$
\left[\begin{array}{c}
E_{x}(\mathbf{r})  \tag{8.33}\\
E_{y}(\mathbf{r}) \\
E_{z}(\mathbf{r})
\end{array}\right] \stackrel{?}{=}\left[\begin{array}{l}
E_{x}\left(\mathcal{R}^{-1}(\mathbf{r})\right) \\
E_{y}\left(\mathcal{R}^{-1}(\mathbf{r})\right) \\
E_{z}\left(\mathcal{R}^{-1}(\mathbf{r})\right)
\end{array}\right] .
$$

No. Because the components get mixed as well as the positions at which those components are evaluated. We will work with many component wave functions, some of which will behave like vectors, and will have to develop the methods and language to tackle this.

## 8.3 multiple wavefunction spaces.

Reading: See $\S 26.5$ in the text [4].
We identified

$$
\begin{equation*}
\psi(\mathbf{r})=\langle\mathbf{r} \mid \psi\rangle, \tag{8.34}
\end{equation*}
$$

with improper basis kets

$$
\begin{equation*}
|\mathbf{r}\rangle . \tag{8.35}
\end{equation*}
$$

Now introduce many function spaces

$$
\left[\begin{array}{c}
\psi_{1}(\mathbf{r})  \tag{8.3}\\
\psi_{2}(\mathbf{r}) \\
\vdots \\
\psi_{\gamma}(\mathbf{r})
\end{array}\right],
$$

with improper (unnormalizable) basis kets

$$
\begin{equation*}
|\mathbf{r} \alpha\rangle, \quad \alpha \in 1,2, \ldots \gamma \tag{8.37}
\end{equation*}
$$

and form

$$
\begin{equation*}
\psi_{\alpha}(\mathbf{r})=\langle\mathbf{r} \alpha \mid \psi\rangle, \tag{8.38}
\end{equation*}
$$

for an abstract ket $|\psi\rangle$.
We will try taking this Hilbert space

$$
\begin{equation*}
H=H_{o} \otimes H_{s}, \tag{8.39}
\end{equation*}
$$

where $H_{o}$ is the Hilbert space of "scalar" QM, " o " orbital and translational motion, associated with kets $|\mathbf{r}\rangle$ and $H_{s}$ is the Hilbert space associated with the $\gamma$ components $|\alpha\rangle$. This latter space we will label the "spin" or "internal physics" (class suggestion: or perhaps intrinsic). This is "unconnected" with translational motion. We build up the basis kets for $H$ by direct products

$$
\begin{equation*}
|\mathbf{r} \alpha\rangle=|\mathbf{r}\rangle \otimes|\alpha\rangle . \tag{8.40}
\end{equation*}
$$

Now, for a rotated ket we seek a general angular momentum operator $\mathbf{J}$ such that

$$
\begin{equation*}
\left|\psi^{\prime}\right\rangle=e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{J} / \hbar}|\psi\rangle \tag{8.41}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathbf{J}=\mathbf{L}+\mathbf{S} \tag{8.42}
\end{equation*}
$$

where $\mathbf{L}$ acts over kets in $H_{o}$, "orbital angular momentum", and $\mathbf{S}$ is the "spin angular momentum", acting on kets in $H_{s}$. Strictly speaking this would be written as direct products involving the respective identities

$$
\begin{equation*}
\mathbf{J}=\mathbf{L} \otimes I_{S}+I_{o} \otimes \mathbf{S} \tag{8.43}
\end{equation*}
$$

We require

$$
\begin{equation*}
\left[J_{i}, J_{j}\right]=i \hbar \sum \epsilon_{i j k} J_{k} \tag{8.44}
\end{equation*}
$$

since $\mathbf{L}$ and $\mathbf{S}$ "act over separate Hilbert spaces". Since these come from legacy operators

$$
\begin{equation*}
\left[L_{i}, S_{j}\right]=0 \tag{8.45}
\end{equation*}
$$

we also know that

$$
\begin{equation*}
\left[L_{i}, L_{j}\right]=i \hbar \sum \epsilon_{i j k} L_{k} \tag{8.46}
\end{equation*}
$$

So

$$
\begin{equation*}
\left[S_{i}, S_{j}\right]=i \hbar \sum \epsilon_{i j k} S_{k} \tag{8.47}
\end{equation*}
$$

as expected. We could, in principle, have more complicated operators, where this would not be true. This is a proposal of sorts. Given such a definition of operators, let us see where we can go with it.

For matrix elements of $\mathbf{L}$ we have

$$
\begin{equation*}
\langle\mathbf{r}| L_{x}\left|\mathbf{r}^{\prime}\right\rangle=-i \hbar\left(y \frac{\partial}{\partial z}-z \frac{\partial}{\partial y}\right) \delta\left(\mathbf{r}-\mathbf{r}^{\prime}\right) \tag{8.48}
\end{equation*}
$$

What are the matrix elements of $\langle\alpha| S_{i}\left|\alpha^{\prime}\right\rangle$ ? From the commutation relationships we know

$$
\begin{align*}
& \sum_{\alpha^{\prime \prime}=1}^{\gamma}\langle\alpha| S_{i}\left|\alpha^{\prime \prime}\right\rangle\left\langle\alpha^{\prime \prime}\right| S_{j}\left|\alpha^{\prime}\right\rangle-\sum_{\alpha^{\prime \prime}=1}^{\gamma}\langle\alpha| S_{j}\left|\alpha^{\prime \prime}\right\rangle\left\langle\alpha^{\prime \prime}\right| S_{i}\left|\alpha^{\prime}\right\rangle  \tag{8.49}\\
&=i \hbar \sum_{k} \epsilon_{i j k}\langle\alpha| S_{k}\left|\alpha^{\prime \prime}\right\rangle
\end{align*}
$$

We see that our matrix element is tightly constrained by our choice of commutator relationships. We have $\gamma^{2}$ such matrix elements, and it turns out that it is possible to choose (or find) matrix elements that satisfy these constraints?

The $\langle\alpha| S_{i}\left|\alpha^{\prime}\right\rangle$ matrix elements that satisfy these constraints are found by imposing the commutation relations

$$
\begin{equation*}
\left[S_{i}, S_{j}\right]=i \hbar \sum \epsilon_{i j k} S_{k} \tag{8.50}
\end{equation*}
$$

and with

$$
\begin{equation*}
S^{2}=\sum_{j} S_{j}^{2} \tag{8.51}
\end{equation*}
$$

(this is just a definition). We find

$$
\begin{equation*}
\left[S^{2}, S_{i}\right]=0, \tag{8.52}
\end{equation*}
$$

and seek eigenkets

$$
\begin{align*}
S^{2}\left|s m_{s}\right\rangle & =s(s+1) \hbar^{2}\left|s m_{s}\right\rangle  \tag{8.53}\\
S_{z}\left|s m_{s}\right\rangle & =\hbar m_{s}\left|s m_{s}\right\rangle .
\end{align*}
$$

Find solutions for $s=1 / 2,1,3 / 2,2, \cdots$, where $m_{s} \in\{-s, \cdots, s\}$. ie. $2 s+1$ possible vectors $\left|s m_{s}\right\rangle$ for a given $s$.

$$
\begin{align*}
& s=\frac{1}{2} \Longrightarrow \gamma=2 \\
& s=1 \Longrightarrow \gamma=3  \tag{8.54}\\
& s=\frac{3}{2} \Longrightarrow \gamma=4 .
\end{align*}
$$

We start with the algebra (mathematically the Lie algebra), and one can compute the Hilbert spaces that are consistent with these algebraic constraints.

We assume that for any type of given particle $S$ is fixed, where this has to do with the nature of the particle

$$
\begin{align*}
s=\frac{1}{2} & \text { A spin 1/2 particle } \\
s=1 & \text { A spin 1 particle }  \tag{8.55}\\
s=\frac{3}{2} & \text { A spin 3/2 particle. }
\end{align*}
$$

The spin $s$ is fixed once we decide that we are talking about a specific type of particle.

A non-relativistic particle in this framework has two nondynamical quantities. One is the mass $m$ and we now introduce a new invariant, the $\operatorname{spin} s$ of the particle.

This has been introduced as a kind of strategy. It is something that we are going to try, and it turns out that it does. This agrees well with experiment.

In 1939 Wigner asked, "what constraints do I get if I constrain the constraints of quantum mechanics with special relativity." It turns out that in the non-relativistic limit, we get just this.

There is a subtlety here, because we get into some logical trouble with the photon with a rest mass of zero ( $m=0$ is certainly allowed as a value of our invariant $m$ above). We can not stop or slow down a photon, so orbital angular momentum is only a conceptual idea. Really, the orbital angular momentum and the spin angular momentum cannot be separated out for a photon, so talking of a spin 1 particle really means spin as in $\mathbf{J}$, and not spin as in $\mathbf{L}$.

Spin one half particles. Reading: See §26.6 in the text [4].
Let us start talking about the simplest case. This includes electrons, all leptons (integer spin particles like photons and the weakly interacting W and Z bosons), and quarks. That is

$$
\begin{align*}
s & =\frac{1}{2} \\
m_{s} & = \pm \frac{1}{2} \tag{8.56}
\end{align*}
$$

with states

$$
\begin{equation*}
\left|s m_{s}\right\rangle=\left|\frac{1}{2}, \frac{1}{2}\right\rangle,\left|\frac{1}{2},-\frac{1}{2}\right\rangle . \tag{8.57}
\end{equation*}
$$

We utilize the convention

$$
\begin{align*}
& \left|\frac{1}{2} \frac{1}{2}\right\rangle=\left|\frac{1}{2},-\frac{1}{2}\right\rangle  \tag{8.58}\\
& \left|\frac{1}{2} \frac{1}{2}\right\rangle=\left|\frac{1}{2}, \frac{1}{2}\right\rangle .
\end{align*}
$$

We calculate

$$
\begin{align*}
S^{2}\left|\frac{1}{2} m_{s}\right\rangle & =\frac{1}{2}\left(\frac{1}{2}+1\right) \hbar^{2}\left|\frac{1}{2} m_{s}\right\rangle  \tag{8.59}\\
& =\frac{3}{4} \hbar^{2}\left|\frac{1}{2} m_{s}\right\rangle,
\end{align*}
$$

and

$$
\begin{equation*}
S_{z}\left|\frac{1}{2} m_{s}\right\rangle=m_{s} \hbar\left|\frac{1}{2} m_{s}\right\rangle \tag{8.60}
\end{equation*}
$$

Now introduce additional shorthand

$$
\begin{align*}
& \left|\frac{1}{2} \frac{1}{2}\right\rangle=|+\rangle \\
& \left.\left|\frac{1}{2} \frac{1}{2}\right\rangle\right\rangle=|-\rangle . \tag{8.61}
\end{align*}
$$

The matrix representations of the spin operators was

$$
S^{2} \rightarrow \frac{3}{4} \hbar^{2}\left[\begin{array}{ll}
1 & 0  \tag{8.62}\\
0 & 1
\end{array}\right],
$$

and

$$
S_{z} \rightarrow \frac{\hbar}{2}\left[\begin{array}{cc}
1 & 0  \tag{8.63}\\
0 & -1
\end{array}\right]
$$

One can easily work out from the commutation relationships that the $x$ and $y$ oriented spin operators are

$$
S_{x} \rightarrow \frac{\hbar}{2}\left[\begin{array}{ll}
0 & 1  \tag{8.64}\\
1 & 0
\end{array}\right]
$$

and

$$
S_{y} \rightarrow \frac{\hbar}{2}\left[\begin{array}{cc}
0 & -i  \tag{8.65}\\
i & 0
\end{array}\right] .
$$

Our next topic will be adding $\mathbf{L}$ into the mix.

### 9.1 REPRESENTATION OF KETS.

Reading: §5.1-§5.9 and §26 in [4].
We found the representations of the spin operators

$$
\begin{align*}
S_{x} & \rightarrow \frac{\hbar}{2}\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right]  \tag{9.1}\\
S_{y} & \rightarrow \frac{\hbar}{2}\left[\begin{array}{cc}
0 & -i \\
i & 0
\end{array}\right]  \tag{9.2}\\
S_{z} & \rightarrow \frac{\hbar}{2}\left[\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right] . \tag{9.3}
\end{align*}
$$

How about kets? For example for $|\chi\rangle \in H_{s}$

$$
|\chi\rangle \rightarrow\left[\begin{array}{l}
\langle+\mid \chi\rangle  \tag{9.4}\\
\langle-\mid \chi\rangle
\end{array}\right],
$$

and

So, for example

$$
S_{y}|+\rangle \rightarrow \frac{\hbar}{2}\left[\begin{array}{cc}
0 & -i  \tag{9.7}\\
i & 0
\end{array}\right] \cdot\left[\begin{array}{l}
1 \\
0
\end{array}\right]=\frac{i \hbar}{2}\left[\begin{array}{l}
0 \\
1
\end{array}\right]
$$

Kets in $H_{o} \otimes H_{s}$

$$
|\psi\rangle \rightarrow\left[\begin{array}{l}
\langle\mathbf{r}+\mid \psi\rangle  \tag{9.8}\\
\langle\mathbf{r}-\mid \psi\rangle
\end{array}\right]=\left[\begin{array}{l}
\psi_{+}(\mathbf{r}) \\
\psi_{-}(\mathbf{r})
\end{array}\right]
$$

This is a "spinor" Put

$$
\langle\mathbf{r} \pm \mid \psi\rangle=\psi_{ \pm}(\mathbf{r})=\psi_{+}\left[\begin{array}{l}
1  \tag{9.9}\\
0
\end{array}\right]+\psi_{-}\left[\begin{array}{l}
0 \\
1
\end{array}\right]
$$

where we use a unit normalization

$$
\begin{equation*}
\langle\psi \mid \psi\rangle=1 . \tag{9.10}
\end{equation*}
$$

Use

$$
\begin{align*}
I & =I_{o} \otimes I_{s} \\
& =\int d^{3} \mathbf{r}|\mathbf{r}\rangle\langle\mathbf{r}| \otimes(|+\rangle\langle+|+|-\rangle\langle-|) \\
& =\int d^{3} \mathbf{r}|\mathbf{r}\rangle\langle\mathbf{r}| \otimes \sum_{\sigma= \pm}|\sigma\rangle\langle\sigma|  \tag{9.11}\\
& =\sum_{\sigma= \pm} \int d^{3} \mathbf{r}|\mathbf{r} \sigma\rangle\langle\mathbf{r} \sigma|,
\end{align*}
$$

so

$$
\begin{align*}
\langle\psi| I|\psi\rangle & =\sum_{\sigma= \pm} \int d^{3} \mathbf{r}\langle\psi \mid \mathbf{r} \sigma\rangle\langle\mathbf{r} \sigma \mid \psi\rangle \\
& =\int d^{3} \mathbf{r}\left(\left|\psi_{+}(\mathbf{r})\right|^{2}+\left|\psi_{-}(\mathbf{r})\right|^{2}\right) . \tag{9.12}
\end{align*}
$$

Alternatively.

$$
\begin{align*}
|\psi\rangle & =I|\psi\rangle \\
& =\int d^{3} \mathbf{r} \sum_{\sigma= \pm}|\mathbf{r} \sigma\rangle\langle\mathbf{r} \sigma \mid \psi\rangle \\
& =\sum_{\sigma= \pm}\left(\int d^{3} \mathbf{r} \psi_{\sigma}(\mathbf{r})\right)|\mathbf{r} \sigma\rangle  \tag{9.13}\\
& =\sum_{\sigma= \pm}\left(\int d^{3} \mathbf{r} \psi_{\sigma}(\mathbf{r})|\mathbf{r}\rangle\right) \otimes|\sigma\rangle .
\end{align*}
$$

In braces we have a ket in $H_{o}$, let us call it

$$
\begin{equation*}
\left|\psi_{\sigma}\right\rangle=\int d^{3} \mathbf{r} \psi_{\sigma}(\mathbf{r})|\mathbf{r}\rangle, \tag{9.14}
\end{equation*}
$$

then

$$
\begin{equation*}
|\psi\rangle=\left|\psi_{+}\right\rangle|+\rangle+\left|\psi_{-}\right\rangle|-\rangle, \tag{9.15}
\end{equation*}
$$

where the direct product $\otimes$ is implied.
We can form a ket in $H_{s}$ as

$$
\begin{equation*}
\langle\mathbf{r} \mid \psi\rangle=\psi_{+}(\mathbf{r})|+\rangle+\psi_{-}(\mathbf{r})|-\rangle . \tag{9.16}
\end{equation*}
$$

An operator $O_{o}$ which acts on $H_{o}$ alone can be promoted to $O_{o} \otimes I_{s}$, which is now an operator that acts on $H_{o} \otimes H_{s}$. We are sometimes a little cavalier in notation and leave this off, but we should remember this. Acting with $O_{o}$ on our ket, we have

$$
\begin{equation*}
O_{o}|\psi\rangle=\left(O_{o}|\psi+\rangle\right)|+\rangle-\left(O_{o}|\psi-\rangle\right)|-\rangle, \tag{9.17}
\end{equation*}
$$

and likewise

$$
\begin{equation*}
O_{s}|\psi\rangle=|\psi+\rangle\left(O_{s}|+\rangle\right)+|\psi-\rangle\left(O_{s}|-\rangle\right), \tag{9.18}
\end{equation*}
$$

so the composite action is

$$
\begin{equation*}
O_{o} O_{s}|\psi\rangle=\left(O_{o}|\psi+\rangle\right)\left(O_{s}|+\rangle\right)+\left(O_{o}|\psi-\rangle\right)\left(O_{s}|-\rangle\right) . \tag{9.19}
\end{equation*}
$$

Suppose we want to rotate a ket, we do this with a full angular momentum operator

$$
\begin{equation*}
e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{J} / \hbar}|\psi\rangle=e^{-i \theta \hat{\mathbf{n}} \mathbf{L} / \hbar e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{S} / \hbar}|\psi\rangle, ~} \tag{9.20}
\end{equation*}
$$

(recalling that $\mathbf{L}$ and $\mathbf{S}$ commute.) So
$e^{-i \theta \hat{\mathbf{n}} \mathbf{J} / \hbar}|\psi\rangle=\left(e^{-i \theta \hat{\mathbf{n}} \mathbf{L} / \hbar}|\psi+\rangle\right)\left(e^{-i \theta \hat{n} \mathbf{S} / \hbar}|+\rangle\right)+\left(e^{-i \theta \hat{\mathbf{n}} \mathbf{L} / \hbar}|\psi-\rangle\right)\left(e^{-i \theta \hat{\mathbf{n}} \mathbf{S} / \hbar}|-\rangle\right)$.

A simple example. Let

$$
\begin{equation*}
|\psi\rangle=\left|\psi_{+}\right\rangle|+\rangle+\left|\psi_{-}\right\rangle|-\rangle . \tag{9.22}
\end{equation*}
$$

Suppose

$$
\begin{align*}
& \left|\psi_{+}\right\rangle=\alpha\left|\psi_{0}\right\rangle  \tag{9.23}\\
& \left|\psi_{-}\right\rangle=\beta\left|\psi_{0}\right\rangle, \tag{9.24}
\end{align*}
$$

where

$$
\begin{equation*}
|\alpha|^{2}+|\beta|^{2}=1 . \tag{9.25}
\end{equation*}
$$

Now let

$$
\begin{equation*}
|\psi\rangle=\left|\psi_{0}\right\rangle|\chi\rangle, \tag{9.26}
\end{equation*}
$$

where

$$
\begin{equation*}
|\chi\rangle=\alpha|+\rangle+\beta|-\rangle \tag{9.27}
\end{equation*}
$$

Imposition of the normalization condition

$$
\begin{equation*}
\langle\psi \mid \psi\rangle=1 \tag{9.28}
\end{equation*}
$$

implies

$$
\begin{equation*}
\left\langle\psi_{0} \mid \psi_{0}\right\rangle\langle\chi \mid \chi\rangle=1 \tag{9.29}
\end{equation*}
$$

so

$$
\begin{equation*}
\left\langle\psi_{0} \mid \psi_{0}\right\rangle=1 \tag{9.30}
\end{equation*}
$$

We are going to concentrate on the unentangled state of eq. (9.26).

- How about with

$$
\begin{equation*}
|\alpha|^{2}=1, \beta=0 \tag{9.31}
\end{equation*}
$$

$|\chi\rangle$ is an eigenket of $S_{z}$ with eigenvalue $\hbar / 2$.

$$
\begin{equation*}
|\beta|^{2}=1, \alpha=0 \tag{9.32}
\end{equation*}
$$

$|\chi\rangle$ is an eigenket of $S_{z}$ with eigenvalue $-\hbar / 2$.

- What is $|\chi\rangle$ if it is an eigenket of $\hat{\mathbf{n}} \cdot \mathbf{S}$ ?

FIXME: F1: standard spherical projection picture, with $\hat{\mathbf{n}}$ projected down onto the $x, y$ plane at angle $\phi$ and at an angle $\theta$ from the $z$ axis.

The eigenvalues will still be $\pm \hbar / 2$ since there is nothing special about the $z$ direction.

$$
\begin{align*}
\hat{\mathbf{n}} \cdot \mathbf{S} & =n_{x} S_{x}+n_{y} S_{y}+n_{z} S_{z} \\
& \rightarrow \frac{\hbar}{2}\left[\begin{array}{cc}
n_{z} & n_{x}-i n_{y} \\
n_{x}+i n_{y} & -n_{z}
\end{array}\right]  \tag{9.33}\\
& =\frac{\hbar}{2}\left[\begin{array}{cc}
\cos \theta & \sin \theta e^{-i \phi} \\
\sin \theta e^{i \phi} & -\cos \theta
\end{array}\right] .
\end{align*}
$$

To find the eigenkets we diagonalize this, and we find representations of the eigenkets are

$$
\begin{align*}
& |\hat{\mathbf{n}}+\rangle \rightarrow\left[\begin{array}{c}
\cos \left(\frac{\theta}{2}\right) e^{-i \phi / 2} \\
\sin \left(\frac{\theta}{2}\right) e^{i \phi / 2}
\end{array}\right]  \tag{9.34}\\
& |\hat{\mathbf{n}}-\rangle \rightarrow\left[\begin{array}{c}
-\sin \left(\frac{\theta}{2}\right) e^{-i \phi / 2} \\
\cos \left(\frac{\theta}{2}\right) e^{i \phi / 2}
\end{array}\right], \tag{9.35}
\end{align*}
$$

with eigenvalues $\hbar / 2$ and $-\hbar / 2$ respectively.
Tossing the specific representation, we have the following abstract representation of states

$$
\begin{align*}
|\hat{\mathbf{n}}+\rangle & \rightarrow \cos \left(\frac{\theta}{2}\right) e^{-i \phi / 2}|+\rangle \sin \left(\frac{\theta}{2}\right) e^{i \phi / 2}|-\rangle  \tag{9.36}\\
|\hat{\mathbf{n}}-\rangle & \rightarrow-\sin \left(\frac{\theta}{2}\right) e^{-i \phi / 2}|+\rangle \cos \left(\frac{\theta}{2}\right) e^{i \phi / 2}|-\rangle \tag{9.37}
\end{align*}
$$

### 9.2 REPRESENTATION OF TWO STATE KETS.

Every ket $|\chi\rangle \rightarrow\left[\begin{array}{l}\alpha \\ \beta\end{array}\right]$ for which $|\alpha|^{2}+|\beta|^{2}=1$ can be written in the form eq. (9.34) for some $\theta$ and $\phi$, neglecting an overall phase factor.

For any ket in $H_{s}$, that ket is "spin up" in some direction.
FIXME: show this.

### 9.3 PAULI SPIN MATRICES.

It is useful to write

$$
\begin{align*}
S_{x} & =\frac{\hbar}{2}\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right] \equiv \frac{\hbar}{2} \sigma_{x}  \tag{9.38}\\
S_{y} & =\frac{\hbar}{2}\left[\begin{array}{cc}
0 & -i \\
i & 0
\end{array}\right] \equiv \frac{\hbar}{2} \sigma_{y}  \tag{9.39}\\
& =\frac{\hbar}{2}\left[\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right] \equiv \frac{\hbar}{2} \sigma_{z}, \tag{9.40}
\end{align*}
$$

where

$$
\begin{align*}
& \sigma_{x}=\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right]  \tag{9.41}\\
& \sigma_{y}=\left[\begin{array}{cc}
0 & -i \\
i & 0
\end{array}\right]  \tag{9.42}\\
& \sigma_{z}=\left[\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right] . \tag{9.43}
\end{align*}
$$

These are the Pauli spin matrices.

Interesting properties.

- Commutators are zero

$$
\begin{equation*}
\left[\sigma_{i}, \sigma_{j}\right]=\sigma_{i} \sigma_{j}+\sigma_{j} \sigma_{i}=0, \quad \text { if } i<j \tag{9.44}
\end{equation*}
$$

- Relations between products have the form

$$
\begin{equation*}
\sigma_{x} \sigma_{y}=i \sigma_{z} \tag{9.45}
\end{equation*}
$$

where the other products can be found by cyclic permutations of the indexes.

- All the Pauli matrices are traceless

$$
\begin{equation*}
\operatorname{tr}\left(\sigma_{i}\right)=0 \tag{9.46}
\end{equation*}
$$

- Pauli "vectors" square to the identity matrix

$$
\begin{equation*}
(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma})^{2}=\sigma_{0} \tag{9.47}
\end{equation*}
$$

where

$$
\begin{equation*}
\hat{\mathbf{n}} \cdot \boldsymbol{\sigma} \equiv n_{x} \sigma_{x}+n_{y} \sigma_{y}+n_{z} \sigma_{z} \tag{9.48}
\end{equation*}
$$

and

$$
\sigma_{0}=\left[\begin{array}{ll}
1 & 0  \tag{9.49}\\
0 & 1
\end{array}\right]
$$

Note that $\operatorname{tr}\left(\sigma_{0}\right) \neq 0$.

- Commutators are

$$
\begin{align*}
{\left[\sigma_{i}, \sigma_{j}\right] } & =2 \delta_{i j} \sigma_{0}  \tag{9.50}\\
{\left[\sigma_{x}, \sigma_{y}\right] } & =2 i \sigma_{z} \tag{9.51}
\end{align*}
$$

where we may also cyclically permute in the indexes of the latter. We can combine these to show that

$$
\begin{equation*}
(\mathbf{A} \cdot \boldsymbol{\sigma})(\mathbf{B} \cdot \boldsymbol{\sigma})=(\mathbf{A} \cdot \mathbf{B}) \sigma_{0}+i(\mathbf{A} \times \mathbf{B}) \cdot \boldsymbol{\sigma} \tag{9.52}
\end{equation*}
$$

where $\mathbf{A}$ and $\mathbf{B}$ are vectors (or more generally operators that commute with the $\sigma$ matrices).

- Product traces are given by

$$
\begin{equation*}
\operatorname{tr}\left(\sigma_{i} \sigma_{j}\right)=2 \delta_{i j} \tag{9.53}
\end{equation*}
$$

- The product trace relation above may be extended to $\sigma_{0}$

$$
\begin{equation*}
\operatorname{tr}\left(\sigma_{\alpha} \sigma_{\beta}\right)=2 \delta_{\alpha \beta} \tag{9.54}
\end{equation*}
$$

where $\alpha, \beta=0, x, y, z$.
Note that any complex matrix $M$ can be written as

$$
\begin{align*}
M & =\sum_{\alpha} m_{a} \sigma_{\alpha} \\
& =\left[\begin{array}{ll}
m_{0}+m_{z} & m_{x}-i m_{y} \\
m_{x}+i m_{y} & m_{0}-m_{z}
\end{array}\right] \tag{9.55}
\end{align*}
$$

for any four complex numbers $m_{0}, m_{x}, m_{y}, m_{z}$, where

$$
\begin{equation*}
m_{\beta}=\frac{1}{2} \operatorname{tr}\left(M \sigma_{\beta}\right) \tag{9.56}
\end{equation*}
$$

### 10.1 FORMAL TAYLOR SERIES EXPANSION.

READING: $\S 27.5$ in the text [4].
We can formally expand our rotation operator in Taylor series

$$
\begin{equation*}
e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{S} / \hbar}=I+(-i \theta \hat{\mathbf{n}} \cdot \mathbf{S} / \hbar)+\frac{1}{2!}(-i \theta \hat{\mathbf{n}} \cdot \mathbf{S} / \hbar)^{2}+\frac{1}{3!}(-i \theta \hat{\mathbf{n}} \cdot \mathbf{S} / \hbar)^{3}+\cdots \tag{10.1}
\end{equation*}
$$

or

$$
\begin{align*}
e^{-i \theta \hat{\mathbf{n}} \cdot \boldsymbol{\sigma} / 2} & =I+(-i \theta \hat{\mathbf{n}} \cdot \boldsymbol{\sigma} / 2)+\frac{1}{2!}(-i \theta \hat{\mathbf{n}} \cdot \boldsymbol{\sigma} / 2)^{2}+\frac{1}{3!}(-i \theta \hat{\mathbf{n}} \cdot \sigma / 2)^{3}+\cdots \\
& =\sigma_{0}+\left(\frac{-i \theta}{2}\right)(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma})+\frac{1}{2!}\left(\frac{-i \theta}{2}\right)(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma})^{2}+\frac{1}{3!}\left(\frac{-i \theta}{2}\right)(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma})^{3}+\cdots \\
& =\sigma_{0}+\left(\frac{-i \theta}{2}\right)(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma})+\frac{1}{2!}\left(\frac{-i \theta}{2}\right) \sigma_{0}+\frac{1}{3!}\left(\frac{-i \theta}{2}\right)(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma})+\cdots \\
& =\sigma_{0}\left(1-\frac{1}{2!}\left(\frac{\theta}{2}\right)^{2}+\cdots\right)-i(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma})\left(\frac{\theta}{2}-\frac{1}{3!}\left(\frac{\theta}{2}\right)^{3}+\cdots\right) \\
& =\cos (\theta / 2) \sigma_{0}-i \sin (\theta / 2)(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma}) \tag{10.2}
\end{align*}
$$

where we have used the fact that $(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma})^{2}=\sigma_{0}$. The representation of the spin operator is thus

$$
\begin{align*}
e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{S} / \hbar} & \rightarrow \cos (\theta / 2) \sigma_{0}-i \sin (\theta / 2)(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma}) \\
& =\cos (\theta / 2) \sigma_{0}-i \sin (\theta / 2)\left(n_{x}\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right]+n_{y}\left[\begin{array}{cc}
0 & -i \\
i & 0
\end{array}\right]+n_{z}\left[\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right]\right) \\
& =\left[\begin{array}{cc}
\cos (\theta / 2)-i n_{z} \sin (\theta / 2) & -i\left(n_{x}-i n_{y}\right) \sin (\theta / 2) \\
-i\left(n_{x}+i n_{y}\right) \sin (\theta / 2) & \cos (\theta / 2)+i n_{z} \sin (\theta / 2)
\end{array}\right] \tag{10.3}
\end{align*}
$$

In particular,

$$
\begin{equation*}
e^{-2 \pi i \hat{\mathbf{n}} \cdot \mathbf{S} / \hbar} \rightarrow \cos \pi \sigma_{0}=-\sigma_{0} \tag{10.4}
\end{equation*}
$$

This "rotates" the ket, but introduces a phase factor.
This can be generalized to other degrees of spin, for $s=1 / 2,3 / 2,5 / 2, \cdots$.

Unfortunate interjection by me. I mentioned the half angle rotation operator that requires a half angle operator sandwich. Prof. Sipe thought I might be talking about a Heisenberg picture representation, where we have something like this in expectation values

$$
\begin{equation*}
\left|\psi^{\prime}\right\rangle=e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{J} / \hbar}|\psi\rangle \tag{10.5}
\end{equation*}
$$

so that

$$
\begin{equation*}
\left\langle\psi^{\prime}\right| O\left|\psi^{\prime}\right\rangle=\langle\psi| e^{i \theta \hat{\mathbf{n}} \cdot \mathbf{J} / \hbar} O e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{J} / \hbar}|\psi\rangle \tag{10.6}
\end{equation*}
$$

However, what I was referring to, was that a general rotation of a vector in a Pauli matrix basis

$$
\begin{equation*}
R\left(\sum a_{k} \sigma_{k}\right)=R(\mathbf{a} \cdot \sigma) \tag{10.7}
\end{equation*}
$$

can be expressed by sandwiching the Pauli vector representation by two half angle rotation operators like our spin $1 / 2$ operators from class today

$$
\begin{equation*}
R(\mathbf{a} \cdot \boldsymbol{\sigma})=e^{-\theta \hat{\mathbf{u}} \cdot \boldsymbol{\sigma} \cdot \boldsymbol{\hat { v }} \cdot \boldsymbol{\sigma}} \mathbf{a} \cdot \boldsymbol{\sigma} e^{\theta \hat{\mathbf{u}} \cdot \boldsymbol{\sigma} \cdot \hat{\mathbf{v}} \cdot \boldsymbol{\sigma} / 2} \tag{10.8}
\end{equation*}
$$

where $\hat{\mathbf{u}}$ and $\hat{\mathbf{v}}$ are two non-colinear orthogonal unit vectors that define the oriented plane that we are rotating in. For example, rotating in the $x-y$ plane, with $\hat{\mathbf{u}}=\hat{\mathbf{x}}$ and $\hat{\mathbf{v}}=\hat{\mathbf{y}}$, we have

$$
\begin{equation*}
R(\mathbf{a} \cdot \boldsymbol{\sigma})=e^{-\theta \sigma_{1} \sigma_{2} / 2}\left(a_{1} \sigma_{1}+a_{2} \sigma_{2}+a_{3} \sigma_{3}\right) e^{\theta \sigma_{1} \sigma_{2} / 2} \tag{10.9}
\end{equation*}
$$

Observe that these exponentials commute with $\sigma_{3}$, leaving

$$
\begin{align*}
R(\mathbf{a} \cdot \boldsymbol{\sigma}) & =\left(a_{1} \sigma_{1}+a_{2} \sigma_{2}\right) e^{\theta \sigma_{1} \sigma_{2}}+a_{3} \sigma_{3} \\
& =\left(a_{1} \sigma_{1}+a_{2} \sigma_{2}\right)\left(\cos \theta+\sigma_{1} \sigma_{2} \sin \theta\right)+a_{3} \sigma_{3} \\
& =\sigma_{1}\left(a_{1} \cos \theta-a_{2} \sin \theta\right)+\sigma_{2}\left(a_{2} \cos \theta+a_{1} \sin \theta\right)+\sigma_{3}\left(a_{3}\right) \tag{10.10}
\end{align*}
$$

yielding our usual coordinate rotation matrix. Expressed in terms of a unit normal to that plane, we form the normal by multiplication with the unit spatial volume element $I=\sigma_{1} \sigma_{2} \sigma_{3}$. For example:

$$
\begin{equation*}
\sigma_{1} \sigma_{2} \sigma_{3}\left(\sigma_{3}\right)=\sigma_{1} \sigma_{2} \tag{10.11}
\end{equation*}
$$

and can in general write a spatial rotation in a Pauli basis representation as a sandwich of half angle rotation matrix exponentials

$$
\begin{equation*}
R(\mathbf{a} \cdot \boldsymbol{\sigma})=e^{-I \theta(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma}) / 2}(\mathbf{a} \cdot \boldsymbol{\sigma}) e^{I \theta(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma}) / 2} \tag{10.12}
\end{equation*}
$$

when $\hat{\mathbf{n}} \cdot \mathbf{a}=0$ we get the complex-number like single sided exponential rotation exponentials (since $\mathbf{a} \cdot \boldsymbol{\sigma}$ commutes with $\mathbf{n} \cdot \boldsymbol{\sigma}$ in that case)

$$
\begin{equation*}
R(\mathbf{a} \cdot \boldsymbol{\sigma})=(\mathbf{a} \cdot \boldsymbol{\sigma}) e^{I \theta(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma})} \tag{10.13}
\end{equation*}
$$

I believe it was pointed out in one of [7] or [9] that rotations expressed in terms of half angle Pauli matrices has caused some confusion to students of quantum mechanics, because this $2 \pi$ "rotation" only generates half of the full spatial rotation. It was argued that this sort of confusion can be avoided if one observes that these half angle rotations exponentials are exactly what we require for general spatial rotations, and that a pair of half angle operators are required to produce a full spatial rotation.

The book [7] takes this a lot further, and produces a formulation of spin operators that is devoid of the normal scalar imaginary $i$ (using the Clifford algebra spatial unit volume element instead), and also does not assume a specific matrix representation of the spin operators. They argue that this leads to some subtleties associated with interpretation, but at the time I was attempting to read that text I did know enough QM to appreciate what they were doing, and have not had time to attempt a new study of that content.

Asked about this offline, our Professor says, "Yes.... but I think this kind of result is essentially what I was saying about the 'rotation of operators' in lecture. As to 'interpreting' the -1 , there are a number of different strategies and ways of thinking about things. But I think the fact remains that a $2 \pi$ rotation of a spinor replaces the spinor by -1 times itself, no matter how you formulate things."

That this double sided half angle construction to rotate a vector falls out of the Heisenberg picture is interesting. Even in a purely geometric Clifford algebra context, I suppose that a vector can be viewed as an operator (acting on another vector it produces a scalar and a bivector, acting on higher grade algebraic elements one gets $+1,-1$ grade elements as a result). Yet that is something that is true, independent of any quantum mechanics. In the books I mentioned, this was not derived, but instead stated, and then proved. That is something that I think deserves a bit of
exploration. Perhaps there is a more natural derivation possible using infinitesimal arguments ... I had guess that scalar or grade selection would take the place of an expectation value in such a geometric argument.

### 10.2 SPIN DYNAMICs.

At least classically, the angular momentum of charged objects

$$
\begin{equation*}
\boldsymbol{\mu}=I A \mathbf{e}_{\perp}, \tag{10.14}
\end{equation*}
$$

is associated with a magnetic moment as illustrated in fig. 10.1. In our


Figure 10.1: Magnetic moment due to steady state current.
scheme, following the (cgs?) text conventions of [4], where the $\mathbf{E}$ and $\mathbf{B}$ have the same units, we write

$$
\begin{equation*}
\boldsymbol{\mu}=\frac{I A}{c} \mathbf{e}_{\perp} . \tag{10.15}
\end{equation*}
$$

For a charge moving in a circle as in fig. 10.2.


Figure 10.2: Charge moving in circle.

$$
\begin{align*}
I & =\frac{\text { charge }}{\text { time }} \\
& =\frac{\text { distance }}{\text { time }} \frac{\text { charge }}{\text { distance }}  \tag{10.16}\\
& =\frac{q v}{2 \pi r},
\end{align*}
$$

so the magnetic moment is

$$
\begin{align*}
\mu & =\frac{q v}{2 \pi r} \frac{\pi r^{2}}{c} \\
& =\frac{q}{2 m c}(m v r)  \tag{10.17}\\
& =\gamma L .
\end{align*}
$$

Here $\gamma$ is the gyromagnetic ratio. Recall that we have a torque, as shown in fig. 10.3.


Figure 10.3: Induced torque in the presence of a magnetic field.

$$
\begin{equation*}
\mathbf{T}=\boldsymbol{\mu} \times \mathbf{B} \tag{10.18}
\end{equation*}
$$

tending to line up $\boldsymbol{\mu}$ with $\mathbf{B}$. The energy is then

$$
\begin{equation*}
-\boldsymbol{\mu} \cdot \mathbf{B} \tag{10.19}
\end{equation*}
$$

Also recall that this torque leads to precession as shown in fig. 10.4.

$$
\begin{equation*}
\frac{d \mathbf{L}}{d t}=\mathbf{T}=\gamma \mathbf{L} \times \mathbf{B} \tag{10.20}
\end{equation*}
$$

This has precession frequency

$$
\begin{equation*}
\omega=-\gamma \mathbf{B} \tag{10.21}
\end{equation*}
$$

For a current due to a moving electron

$$
\begin{equation*}
\gamma=-\frac{e}{2 m c}<0 \tag{10.22}
\end{equation*}
$$

where we are, here, writing for charge on the electron $-e$.


Figure 10.4: Precession due to torque.

Question: steady state currents only? Yes, this is only true for steady state currents. For the translational motion of an electron, even if it is not moving in a steady way, regardless of its dynamics

$$
\begin{equation*}
\boldsymbol{\mu}_{0}=-\frac{e}{2 m c} \mathbf{L} \tag{10.23}
\end{equation*}
$$

Now, back to quantum mechanics, we turn $\mu_{0}$ into a dipole moment operator and $\mathbf{L}$ is "promoted" to an angular momentum operator

$$
\begin{equation*}
H_{\mathrm{int}}=-\boldsymbol{\mu}_{0} \cdot \mathbf{B} \tag{10.24}
\end{equation*}
$$

What about the "spin"? Perhaps

$$
\begin{equation*}
\boldsymbol{\mu}_{s}=\gamma_{s} \mathbf{S} \tag{10.25}
\end{equation*}
$$

We write this as

$$
\begin{equation*}
\boldsymbol{\mu}_{s}=g\left(-\frac{e}{2 m c}\right) \mathbf{S} \tag{10.26}
\end{equation*}
$$

so that

$$
\begin{equation*}
\gamma_{s}=-\frac{g e}{2 m c} \tag{10.27}
\end{equation*}
$$

Experimentally, one finds to very good approximation

$$
\begin{equation*}
g=2 \tag{10.28}
\end{equation*}
$$

There was a lot of trouble with this in early quantum mechanics where people got things wrong, and cancelled the wrong factors of 2 . In fact, Dirac's relativistic theory for the electron predicts $g=2$. When this is measured experimentally, one does not get exactly $g=2$, and a theory
that also incorporates photon creation and destruction and the interaction with the electron with such (virtual) photons. We get

$$
\begin{align*}
g_{\text {theory }} & =2(1.001159652140( \pm 28))  \tag{10.29}\\
g_{\text {experimental }} & =2(1.0011596521884( \pm 43)) .
\end{align*}
$$

Richard Feynman compared the precision of quantum mechanics, referring to this measurement, "to predicting a distance as great as the width of North America to an accuracy of one human hair's breadth".

## 10.3 the hydrogen atom with spin.

READING: what chapter of [4] ?
For a spinless hydrogen atom, the Hamiltonian was

$$
\begin{equation*}
H=H_{\mathrm{CM}} \otimes H_{\mathrm{rel}}, \tag{10.30}
\end{equation*}
$$

where we have independent Hamiltonian's for the motion of the center of mass and the relative motion of the electron to the proton. The basis kets for these could be designated $\left|\mathbf{p}_{\mathrm{CM}}\right\rangle$ and $\left|\mathbf{p}_{\text {rel }}\right\rangle$ respectively. Now we want to augment this, treating

$$
\begin{equation*}
H=H_{\mathrm{CM}} \otimes H_{\mathrm{rel}} \otimes H_{\mathrm{s}}, \tag{10.31}
\end{equation*}
$$

where $H_{\mathrm{s}}$ is the Hamiltonian for the spin of the electron. We are neglecting the spin of the proton, but that could also be included (this turns out to be a lesser effect). We will introduce a Hamiltonian including the dynamics of the relative motion and the electron spin

$$
\begin{equation*}
H_{\mathrm{rel}} \otimes H_{\mathrm{s}} . \tag{10.32}
\end{equation*}
$$

Covering the Hilbert space for this system we will use basis kets

$$
\begin{align*}
& |n l m \pm\rangle,  \tag{10.33}\\
& |n l m+\rangle \rightarrow\left[\begin{array}{l}
\langle\mathbf{r}+\mid n l m+\rangle \\
\langle\mathbf{r}-\mid n l m+\rangle
\end{array}\right]=\left[\begin{array}{c}
\Phi_{n l m}(\mathbf{r}) \\
0
\end{array}\right] .  \tag{10.34}\\
& |n l m-\rangle \rightarrow\left[\begin{array}{c}
\langle\mathbf{r}+\mid n l m-\rangle \\
\langle\mathbf{r}-\mid n l m-\rangle
\end{array}\right]=\left[\begin{array}{c}
0 \\
\Phi_{n l m}(\mathbf{r})
\end{array}\right] .
\end{align*}
$$

Here $\mathbf{r}$ should be understood to really mean $\mathbf{r}_{\text {rel }}$. Our full Hamiltonian, after introducing a magnetic perturbation is

$$
\begin{equation*}
H=\frac{P_{\mathrm{CM}}^{2}}{2 M}+\left(\frac{P_{\mathrm{rel}}^{2}}{2 \mu}-\frac{e^{2}}{R_{\mathrm{rel}}}\right)-\boldsymbol{\mu}_{0} \cdot \mathbf{B}-\boldsymbol{\mu}_{s} \cdot \mathbf{B} \tag{10.35}
\end{equation*}
$$

where

$$
\begin{equation*}
M=m_{\text {proton }}+m_{\text {electron }} \tag{10.36}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{1}{\mu}=\frac{1}{m_{\text {proton }}}+\frac{1}{m_{\text {electron }}} \tag{10.37}
\end{equation*}
$$

For a uniform magnetic field

$$
\begin{align*}
\boldsymbol{\mu}_{0} & =\left(-\frac{e}{2 m c}\right) \mathbf{L}  \tag{10.38}\\
\boldsymbol{\mu}_{s} & =g\left(-\frac{e}{2 m c}\right) \mathbf{S} \tag{10.39}
\end{align*}
$$

We also have higher order terms (higher order multipoles) and relativistic corrections (like spin orbit coupling [20]).

## 1 <br> 1

 CLEBSCH-GORDON.
### 11.1 TWO SPINS.

READING: §28 of [4].

Example: Consider two electrons, 1 in each of 2 quantum dots.

$$
\begin{equation*}
H=H_{1} \otimes H_{2} \tag{11.1}
\end{equation*}
$$

where $H_{1}$ and $H_{2}$ are both spin Hamiltonian's for respective 2D Hilbert spaces. Our complete Hilbert space is thus a 4D space. We will write

Can introduce

$$
\begin{align*}
& \mathbf{S}_{1}=\mathbf{S}_{1}^{(1)} \otimes I^{(2)} \\
& \mathbf{S}_{2}=I^{(1)} \otimes \mathbf{S}_{2}^{(2)} \tag{11.3}
\end{align*}
$$

Here we "promote" each of the individual spin operators to spin operators in the complete Hilbert space. We write

$$
\begin{align*}
& S_{1 z}|++\rangle=\frac{\hbar}{2}|++\rangle \\
& S_{1 z}|+-\rangle=\frac{\hbar}{2}|+-\rangle \tag{11.4}
\end{align*}
$$

Write

$$
\begin{equation*}
\mathbf{S}=\mathbf{S}_{1}+\mathbf{S}_{2} \tag{11.5}
\end{equation*}
$$

for the full spin angular momentum operator. The $z$ component of this operator is

$$
\begin{equation*}
S_{z}=S_{1 z}+S_{2 z} \tag{11.6}
\end{equation*}
$$

$$
\begin{align*}
& S_{z}|++\rangle=\left(S_{1 z}+S_{2 z}\right)|++\rangle=\left(\frac{\hbar}{2}+\frac{\hbar}{2}\right)|++\rangle=\hbar|++\rangle \\
& S_{z}|+-\rangle=\left(S_{1 z}+S_{2 z}\right)|+-\rangle=\left(\frac{\hbar}{2}-\frac{\hbar}{2}\right)|+-\rangle=0 \\
& S_{z}|-+\rangle=\left(S_{1 z}+S_{2 z}\right)|-+\rangle=\left(-\frac{\hbar}{2}+\frac{\hbar}{2}\right)|-+\rangle=0  \tag{11.7}\\
& S_{z}|--\rangle=\left(S_{1 z}+S_{2 z}\right)|--\rangle=\left(-\frac{\hbar}{2}-\frac{\hbar}{2}\right)|--\rangle=-\hbar|--\rangle .
\end{align*}
$$

So, we find that $|x x\rangle$ are all eigenkets of $S_{z}$. These will also all be eigenkets of $\mathbf{S}_{1}^{2}=S_{1 x}^{2}+S_{1 y}^{2}+S_{1 z}^{2}$ since we have

$$
\begin{align*}
& S_{1}^{2}|x x\rangle=\hbar^{2}\left(\frac{1}{2}\right)\left(1+\frac{1}{2}\right)|x x\rangle=\frac{3}{4} \hbar^{2}|x x\rangle  \tag{11.8}\\
& S_{2}^{2}|x x\rangle=\hbar^{2}\left(\frac{1}{2}\right)\left(1+\frac{1}{2}\right)|x x\rangle=\frac{3}{4} \hbar^{2}|x x\rangle \\
& \mathbf{S}^{2}=\left(\mathbf{S}_{1}+\mathbf{S}_{2}\right) \cdot\left(\mathbf{S}_{1}+\mathbf{S}_{2}\right) \\
& \quad=S_{1}^{2}+S_{2}^{2}+2 \mathbf{S}_{1} \cdot \mathbf{S}_{2} \tag{11.9}
\end{align*}
$$

Note that we have a commutation assumption here $\left[S_{1 i}, S_{2 i}\right]=0$, since we have written $2 \mathbf{S}_{1} \cdot \mathbf{S}_{2}$ instead of $\sum_{i} S_{1 i} S_{2 i}+S_{2 i} S_{1 i}$. The justification for this appears to be the promotion of the spin operators in eq. (11.3) to operators in the complete Hilbert space, since each of these spin operators acts only on the kets associated with their index.

Are all the product kets also eigenkets of $\mathbf{S}^{2}$ ? Calculate

$$
\begin{align*}
\mathbf{S}^{2}|+-\rangle & =\left(S_{1}^{2}+S_{2}^{2}+2 \mathbf{S}_{1} \cdot \mathbf{S}_{2}\right)|+-\rangle \\
& =\left(\frac{3}{4} \hbar^{2}+\frac{3}{4} \hbar^{2}\right)+2 S_{1 x} S_{2 x}|+-\rangle+2 S_{1 y} S_{2 y}|+-\rangle+2 S_{1 z} S_{2 z}|+-\rangle \tag{11.10}
\end{align*}
$$

For the $z$ mixed terms, we have

$$
\begin{equation*}
2 S_{1 z} S_{2 z}|+-\rangle=2\left(\frac{\hbar}{2}\right)\left(-\frac{\hbar}{2}\right)|+-\rangle \tag{11.11}
\end{equation*}
$$

$$
\begin{equation*}
\mathbf{S}^{2}|+-\rangle=\hbar^{2}|+-\rangle+2 S_{1 x} S_{2 x}|+-\rangle+2 S_{1 y} S_{2 y}|+-\rangle \tag{11.12}
\end{equation*}
$$

Since we have set our spin direction in the z direction with
we have

$$
\begin{align*}
& S_{x}|+\rangle \rightarrow \frac{\hbar}{2}\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right]\left[\begin{array}{l}
1 \\
0
\end{array}\right]=\frac{\hbar}{2}\left[\begin{array}{l}
0 \\
1
\end{array}\right]=\frac{\hbar}{2}|-\rangle \\
& S_{x}|-\rangle \rightarrow \frac{\hbar}{2}\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right]\left[\begin{array}{l}
0 \\
1
\end{array}\right]=\frac{\hbar}{2}\left[\begin{array}{l}
1 \\
0
\end{array}\right]=\frac{\hbar}{2}|+\rangle  \tag{11.14}\\
& S_{y}|+\rangle \rightarrow \frac{\hbar}{2}\left[\begin{array}{ll}
0 & -i \\
i & 0
\end{array}\right]\left[\begin{array}{l}
1 \\
0
\end{array}\right]=\frac{i \hbar}{2}\left[\begin{array}{l}
0 \\
1
\end{array}\right]=\frac{i \hbar}{2}|-\rangle \\
& S_{y}|-\rangle \rightarrow \frac{\hbar}{2}\left[\begin{array}{ll}
0 & -i \\
i & 0
\end{array}\right]\left[\begin{array}{l}
0 \\
1
\end{array}\right]=\frac{-i \hbar}{2}\left[\begin{array}{l}
1 \\
0
\end{array}\right]=-\frac{i \hbar}{2}|+\rangle .
\end{align*}
$$

And are able to arrive at the action of $\mathbf{S}^{2}$ on our mixed composite state

$$
\begin{equation*}
\mathbf{S}^{2}|+-\rangle=\hbar^{2}(|+-\rangle+|-+\rangle) \tag{11.15}
\end{equation*}
$$

For the action on the $|++\rangle$ state we have

$$
\begin{align*}
\mathbf{S}^{2}|++\rangle & =\left(\frac{3}{4} \hbar^{2}+\frac{3}{4} \hbar^{2}\right)|++\rangle+2 \frac{\hbar^{2}}{4}|--\rangle+2 i^{2} \frac{\hbar^{2}}{4}|--\rangle+2\left(\frac{\hbar}{2}\right)\left(\frac{\hbar}{2}\right)|++\rangle \\
& =2 \hbar^{2}|++\rangle \tag{11.16}
\end{align*}
$$

and on the $|--\rangle$ state we have

$$
\begin{align*}
\mathbf{S}^{2}|--\rangle= & \left(\frac{3}{4} \hbar^{2}+\frac{3}{4} \hbar^{2}\right)|--\rangle+2 \frac{(-\hbar)^{2}}{4}|++\rangle \\
& +2 i^{2} \frac{\hbar^{2}}{4}|++\rangle+2\left(-\frac{\hbar}{2}\right)\left(-\frac{\hbar}{2}\right)|--\rangle  \tag{11.17}\\
= & 2 \hbar^{2}|--\rangle
\end{align*}
$$

All of this can be assembled into a tidier matrix form

$$
\mathbf{S}^{2} \rightarrow \hbar^{2}\left[\begin{array}{llll}
2 & 0 & 0 & 0  \tag{11.18}\\
0 & 1 & 1 & 0 \\
0 & 1 & 1 & 0 \\
0 & 0 & 0 & 2
\end{array}\right]
$$

where the matrix is taken with respect to the (ordered) basis

However,

$$
\begin{align*}
{\left[\mathbf{S}^{2}, S_{z}\right] } & =0 \\
{\left[S_{i}, S_{j}\right] } & =i \hbar \sum_{k} \epsilon_{i j k} S_{k} \tag{11.20}
\end{align*}
$$

Also, $\left[\mathbf{S}^{2}, S_{i}\right]=0$. It should be possible to find eigenkets of $\mathbf{S}^{2}$ and $S_{z}$

$$
\begin{align*}
\mathbf{S}^{2}\left|s m_{s}\right\rangle & =s(s+1) \hbar^{2}\left|s m_{s}\right\rangle \\
S_{z}\left|s m_{s}\right\rangle & =\hbar m_{s}\left|s m_{s}\right\rangle \tag{11.21}
\end{align*}
$$

An orthonormal set of eigenkets of $\mathbf{S}^{2}$ and $S_{z}$ is found to be

The first three kets here can be grouped into a triplet in a 3D Hilbert space, whereas the last treated as a singlet in a 1D Hilbert space. Form a grouping

$$
\begin{equation*}
H=H_{1} \otimes H_{2} \tag{11.23}
\end{equation*}
$$

Can write

$$
\begin{equation*}
\frac{1}{2} \otimes \frac{1}{2}=1 \oplus 0 \tag{11.24}
\end{equation*}
$$

where the 1 and 0 here refer to the spin index $s$.

Other examples. Consider, perhaps, the $l=5$ state of the hydrogen atom

$$
\begin{align*}
J_{1}^{2}\left|j_{1} m_{1}\right\rangle & =j_{1}\left(j_{1}+1\right) \hbar^{2}\left|j_{1} m_{1}\right\rangle  \tag{11.25}\\
J_{1 z}\left|j_{1} m_{1}\right\rangle & =\hbar m_{1}\left|j_{1} m_{1}\right\rangle, \\
J_{2}^{2}\left|j_{2} m_{2}\right\rangle & =j_{2}\left(j_{2}+1\right) \hbar^{2}\left|j_{2} m_{2}\right\rangle  \tag{11.26}\\
J_{2 z}\left|j_{2} m_{2}\right\rangle & =\hbar m_{2}\left|j_{2} m_{2}\right\rangle .
\end{align*}
$$

Consider the Hilbert space spanned by $\left|j_{1} m_{1}\right\rangle \otimes\left|j_{2} m_{2}\right\rangle$, a $\left(2 j_{1}+1\right)\left(2 j_{2}+1\right)$ dimensional space. How to find the eigenkets of $J^{2}$ and $J_{z}$ ?

### 11.2 MORE ON TWO SPIN SYSTEMS.

READING: Covering § 26.5 of the text [4].

$$
\begin{equation*}
\frac{1}{2} \otimes \frac{1}{2}=1 \oplus 0 \tag{11.27}
\end{equation*}
$$

where 1 is a triplet state for $s=1$ and 0 the "singlet" state with $s=0$. We want to consider the angular momentum of the entire system

$$
\begin{equation*}
j_{1} \otimes j_{2}=? \tag{11.28}
\end{equation*}
$$

Why bother? Often it is true that

$$
\begin{equation*}
[H, \mathbf{J}]=0 \tag{11.29}
\end{equation*}
$$

so, in that case, the eigenstates of the total angular momentum are also energy eigenstates, so considering the angular momentum problem can help in finding these energy eigenstates.

Rotation operator. Form

$$
\begin{equation*}
e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{J} / \hbar} \tag{11.30}
\end{equation*}
$$

where

$$
\begin{equation*}
\hat{\mathbf{n}} \cdot \mathbf{J}=n_{x} J_{x}+n_{y} J_{y}+n_{z} J_{z} . \tag{11.31}
\end{equation*}
$$

Recall the definitions of the raising or lowering operators

$$
\begin{equation*}
J_{ \pm}=J_{x} \pm i J_{y} \tag{11.32}
\end{equation*}
$$

or

$$
\begin{align*}
J_{x} & =\frac{1}{2}\left(J_{+}+J_{-}\right) \\
J_{y} & =\frac{1}{2 i}\left(J_{+}-J_{-}\right) . \tag{11.33}
\end{align*}
$$

We have

$$
\begin{equation*}
\hat{\mathbf{n}} \cdot \mathbf{J}=n_{x} \frac{1}{2}\left(J_{+}+J_{-}\right)+n_{y} \frac{1}{2 i}\left(J_{+}-J_{-}\right)+n_{z} J_{z}, \tag{11.3}
\end{equation*}
$$

and

$$
\begin{equation*}
J_{ \pm}|j m\rangle=\hbar\left((j \mp m)\left(j \pm m_{1}\right)\right)^{1 / 2}|j, m \pm 1\rangle . \tag{11.35}
\end{equation*}
$$

So

$$
\begin{equation*}
\left\langle j^{\prime} m^{\prime}\right| e^{-i \theta \hat{n} \cdot \mathbf{J} / \hbar}|j m\rangle=0, \tag{11.36}
\end{equation*}
$$

unless $j=j^{\prime}$.

$$
\begin{equation*}
\left\langle j m^{\prime}\right| e^{-i \theta \hat{n} \cdot \mathbf{J} / \hbar}|j m\rangle, \tag{11.37}
\end{equation*}
$$

is a $(2 j+1) \times(2 j+1)$ matrix. Combining rotations

$$
\begin{align*}
& \left\langle j m^{\prime}\right| e^{-i \theta_{b} \hat{\mathbf{n}}_{a} \cdot \mathbf{J} / \hbar} e^{-i \theta_{a} \hat{\mathbf{n}}_{b} \cdot \mathbf{J} / \hbar}|j m\rangle \\
& \quad=\sum_{m^{\prime \prime}}\left\langle j m^{\prime}\right| e^{-i \theta_{b} \hat{\mathbf{n}}_{a} \cdot \mathbf{J} / \hbar}\left|j m^{\prime \prime}\right\rangle\left\langle j m^{\prime \prime}\right| e^{-i \theta_{a} \hat{\mathbf{n}}_{b} \cdot \mathbf{J} / \hbar}|j m\rangle . \tag{11.38}
\end{align*}
$$

If

$$
\begin{equation*}
e^{-i \theta \hat{n} \mathbf{U} / \hbar}=e^{-i \theta_{b} \hat{\mathbf{n}}_{a} \cdot \mathbf{J} / \hbar} e^{-i \theta_{a} \hat{\mathbf{n}}_{b} \cdot \mathbf{J} / \hbar}, \tag{11.39}
\end{equation*}
$$

(something that may be hard to compute but possible), then

$$
\begin{equation*}
\left\langle j m^{\prime}\right| e^{-i \theta \hat{n} \mathbf{J} / \hbar}|j m\rangle=\sum_{m^{\prime \prime}}\left\langle j m^{\prime}\right| e^{-i \theta_{b} \hat{\mathbf{n}}_{a} \cdot \mathbf{J} / \hbar}\left|j m^{\prime \prime}\right\rangle\left\langle j m^{\prime \prime}\right| e^{-i \theta_{a} \hat{\mathbf{n}}_{b} \cdot \mathbf{J} / \hbar}|j m\rangle . \tag{11.40}
\end{equation*}
$$

For fixed $j$, the matrices $\left\langle j m^{\prime}\right| e^{-i \theta \hat{n} \cdot \mathbf{J} / \hbar}|j m\rangle$ form a representation of the rotation group. The $(2 j+1)$ representations are irreducible. (This will not be proven).

It may be that there may be big blocks of zeros in some of the matrices, but they cannot be simplified any further?

Back to the two particle system

$$
\begin{equation*}
j_{1} \otimes j_{2}=? \tag{11.41}
\end{equation*}
$$

If we use

$$
\begin{equation*}
\left|j_{1} m_{1}\right\rangle \otimes\left|j_{2} m_{2}\right\rangle \tag{11.42}
\end{equation*}
$$

If a $j_{1}$ and a $j_{2}$ are picked then

$$
\begin{equation*}
\left\langle j_{1} m_{1}^{\prime} ; j_{2} m_{2}^{\prime}\right| e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{J} / \hbar}\left|j_{1} m_{1} ; j_{2} m_{2}\right\rangle \tag{11.43}
\end{equation*}
$$

is also a representation of the rotation group, but these sort of matrices can be simplified a lot. This basis of dimensionality $\left(2 j_{1}+1\right)\left(2 j_{2}+1\right)$ is reducible.

A lot of this is motivation, and we still want a representation of $j_{1} \otimes j_{2}$. Recall that

$$
\begin{equation*}
\frac{1}{2} \otimes \frac{1}{2}=1 \oplus 0=\left(\frac{1}{2}+\frac{1}{2}\right) \oplus\left(\frac{1}{2}-\frac{1}{2}\right) \tag{11.44}
\end{equation*}
$$

Might guess that, for $j_{1} \geq j_{2}$

$$
\begin{equation*}
j_{1} \otimes j_{2}=\left(j_{1}+j_{2}\right) \oplus\left(j_{1}+j_{2}-1\right) \oplus \cdots\left(j_{1}-j_{2}\right) . \tag{11.45}
\end{equation*}
$$

Suppose that this is right. Then

$$
\begin{equation*}
5 \otimes \frac{1}{2}=\frac{11}{2} \oplus \frac{9}{2} \tag{11.46}
\end{equation*}
$$

Check for dimensions:

$$
\begin{align*}
& 1 \otimes 1=2 \oplus 1 \oplus 0  \tag{11.47}\\
& 3 \times 3=5+3+1
\end{align*}
$$

$Q: \quad$ What was this $\oplus$ ? It was just made up. We are creating a shorthand to say that we have a number of different basis states for each of the groupings. Check for dimensions in general

$$
\begin{equation*}
\left(2 j_{1}+1\right)\left(2 j_{2}+1\right) \stackrel{?}{=} \tag{11.48}
\end{equation*}
$$

We find

$$
\begin{align*}
\sum_{j_{1}-j_{2}}^{j_{1}+j_{2}}(2 j+1) & =\sum_{j=0}^{j_{1}+j_{2}}(2 j+1)-\sum_{j=0}^{j_{1}-j_{2}-1}(2 j+1)  \tag{11.49}\\
& =\left(2 j_{1}+1\right)\left(2 j_{2}+1\right)
\end{align*}
$$

Using

$$
\begin{equation*}
\sum_{n=0}^{N} n=\frac{N(N+1)}{2} \tag{11.50}
\end{equation*}
$$

then

$$
\begin{equation*}
j_{1} \otimes j_{2}=\left(j_{1}+j_{2}\right) \oplus\left(j_{1}+j_{2}-1\right) \oplus \cdots\left(j_{1}-j_{2}\right) \tag{11.51}
\end{equation*}
$$

In fact, this is correct. Proof "by construction" to follow

$$
\begin{align*}
\left|j_{1} m_{1}\right\rangle & \otimes\left|j_{2} m_{2}\right\rangle  \tag{11.52}\\
J^{2}|j m\rangle & =j(j+1) \hbar^{2}|j m\rangle  \tag{11.53}\\
J_{z}|j m\rangle & =m \hbar|j m\rangle
\end{align*}
$$

denote also by

$$
\begin{equation*}
\left|j m ; j_{1} j_{2}\right\rangle \tag{11.54}
\end{equation*}
$$

but will often omit the ; $j_{1} j_{2}$ portion. Tab. 11.1 tabulates the $j=j_{1}+j_{2}-k$ progression of kets with $k=0, \cdots 2 j_{2}$

| Table 11.1: Clebsch-Gordan kets. |  |  |  |
| :---: | :---: | :---: | :---: |
| $j=j_{1}+j_{2}$ | $j=j_{1}+j_{2}-1$ | $\ldots$ | $j=j_{1}-j_{2}$ |
| $\left\|j_{1}+j_{2}, j_{1}+j_{2}\right\rangle$ |  |  |  |
| $\left\|j_{1}+j_{2}, j_{1}+j_{2}-1\right\rangle$ | $\left\|j_{1}+j_{2}-1, j_{1}+j_{2}-1\right\rangle$ |  |  |
|  |  |  | $\left\|j_{1}+j_{2}-1, j_{1}+j_{2}-2\right\rangle$ |
| $\vdots$ | $\vdots$ |  | $\left\|j_{1}-j_{2}, j_{1}-j_{2}\right\rangle$ |
| $\vdots$ |  |  | $\vdots$ |
| $\vdots$ |  |  | $\left\|j_{1}-j_{2},-\left(j_{1}-j_{2}\right)\right\rangle$ |
|  |  |  |  |
| $\left\|j_{1}+j_{2},-\left(j_{1}+j_{2}-1\right)\right\rangle$ | $\left\|j_{1}+j_{2}-1,-\left(j_{1}+j_{2}-1\right)\right\rangle$ |  |  |
| $\left\|j_{1}+j_{2},-\left(j_{1}+j_{2}\right)\right\rangle$ |  |  |  |

Look at

$$
\begin{equation*}
\left|j_{1}+j_{2}, j_{1}+j_{2}\right\rangle \tag{11.55}
\end{equation*}
$$

$$
\begin{equation*}
J_{z}\left|j_{1}+j_{2}, j_{1}+j_{2}\right\rangle=\left(j_{1}+j_{2}\right) \hbar\left|j_{1}+j_{2}, j_{1}+j_{2}\right\rangle \tag{11.56}
\end{equation*}
$$

$$
\begin{equation*}
J_{z}\left(\left|j_{1} m_{1}\right\rangle \otimes\left|j_{2} m_{2}\right\rangle\right)=\left(m_{1}+m_{2}\right) \hbar\left(\left|j_{1} m_{1}\right\rangle \otimes\left|j_{2} m_{2}\right\rangle\right) \tag{11.57}
\end{equation*}
$$

we must have

$$
\begin{equation*}
\left|j_{1}+j_{2}, j_{1}+j_{2}\right\rangle=e^{i \phi}\left(\left|j_{1} j_{1}\right\rangle \otimes\left|j_{2} j_{2}\right\rangle\right) \tag{11.58}
\end{equation*}
$$

So $\left|j_{1}+j_{2}, j_{1}+j_{2}\right\rangle$ must be a superposition of states $\left|j_{1} m_{1}\right\rangle \otimes\left|j_{2} m_{2}\right\rangle$ with $m_{1}+m_{2}=j_{1}+j_{2}$. Choosing $e^{i \phi}=1$ is called the Clebsch-Gordan convention.

$$
\begin{equation*}
\left|j_{1}+j_{2}, j_{1}+j_{2}\right\rangle=\left|j_{1} j_{1}\right\rangle \otimes\left|j_{2} j_{2}\right\rangle \tag{11.59}
\end{equation*}
$$

We now move down column.

$$
\begin{equation*}
J_{-}\left|j_{1}+j_{2}, j_{1}+j_{2}\right\rangle=\hbar\left(2\left(j_{1}+j_{2}\right)\right)^{1 / 2}\left|j_{1}+j_{2}, j_{1}+j_{2}-1\right\rangle \tag{11.60}
\end{equation*}
$$

So

$$
\begin{align*}
\left|j_{1}+j_{2}, j_{1}+j_{2}-1\right\rangle & =\frac{J_{-}\left|j_{1}+j_{2}, j_{1}+j_{2}\right\rangle}{\hbar\left(2\left(j_{1}+j_{2}\right)\right)^{1 / 2}}  \tag{11.61}\\
& =\frac{\left(J_{1-}+J_{2-}\right)\left|j_{1} j_{1}\right\rangle \otimes\left|j_{2} j_{2}\right\rangle}{\hbar\left(2\left(j_{1}+j_{2}\right)\right)^{1 / 2}}
\end{align*}
$$

### 11.3 RECAP: TABLE OF TWO SPIN ANGULAR MOMENTA.

Let's review our table above.

First column. Let us start with computation of the kets in the lowest position of the first column, which we will obtain by successive application of the lowering operator to the state

$$
\begin{equation*}
\left|j_{1}+j_{2}, j_{1}+j_{2}\right\rangle=\left|j_{1} j_{1}\right\rangle \otimes\left|j_{2} j_{2}\right\rangle \tag{11.62}
\end{equation*}
$$

Recall that our lowering operator was found to be (or defined as)

$$
\begin{equation*}
J_{-}|j, m\rangle=\sqrt{(j+m)(j-m+1)} \hbar|j, m-1\rangle \tag{11.63}
\end{equation*}
$$

so that application of the lowering operator gives us

$$
\begin{align*}
\mid j_{1}+ & \left.j_{2}, j_{1}+j_{2}-1\right\rangle \\
= & \frac{J_{-}\left|j_{1} j_{1}\right\rangle \otimes\left|j_{2} j_{2}\right\rangle}{\left(2\left(j_{1}+j_{2}\right)\right)^{1 / 2} \hbar} \\
= & \frac{\left(J_{1-}+J_{2-}\right)\left|j_{1} j_{1}\right\rangle \otimes\left|j_{2} j_{2}\right\rangle}{\left(2\left(j_{1}+j_{2}\right)\right)^{1 / 2} \hbar} \\
= & \frac{\left(\sqrt{\left(j_{1}+j_{1}\right)\left(j_{1}-j_{1}+1\right)} \hbar\left|j_{1}\left(j_{1}-1\right)\right\rangle\right) \otimes\left|j_{2} j_{2}\right\rangle}{\left(2\left(j_{1}+j_{2}\right)\right)^{1 / 2} \hbar} \\
& +\frac{\left|j_{1} j_{1}\right\rangle \otimes\left(\sqrt{\left(j_{2}+j_{2}\right)\left(j_{2}-j_{2}+1\right)} \hbar\left|j_{2}\left(j_{2}-1\right)\right\rangle\right)}{\left(2\left(j_{1}+j_{2}\right)\right)^{1 / 2} \hbar} \\
= & \left(\frac{j_{1}}{j_{1}+j_{2}}\right)^{1 / 2}\left|j_{1}\left(j_{1}-1\right)\right\rangle \otimes\left|j_{2} j_{2}\right\rangle+\left(\frac{j_{2}}{j_{1}+j_{2}}\right)^{1 / 2}\left|j_{1} j_{1}\right\rangle \otimes\left|j_{2}\left(j_{2}-1\right)\right\rangle . \tag{11.64}
\end{align*}
$$

Proceeding iteratively would allow us to finish off this column.

Second column. Moving on to the second column, the top most element in the table

$$
\begin{equation*}
\left|j_{1}+j_{2}-1, j_{1}+j_{2}-1\right\rangle \tag{11.65}
\end{equation*}
$$

can only be made up of $\left|j_{1} m_{1}\right\rangle \otimes\left|j_{2} m_{2}\right\rangle$ with $m_{1}+m_{2}=j_{1}+j_{2}-1$. There are two possibilities

$$
\begin{array}{lll}
m_{1}=j_{1} & m_{2}=j_{2}-1  \tag{11.66}\\
m_{1}=j_{1}-1 & m_{2}=j_{2}
\end{array}
$$

So for some $A$ and $B$ to be determined we must have

$$
\begin{equation*}
\left|j_{1}+j_{2}-1, j_{1}+j_{2}-1\right\rangle=A\left|j_{1} j_{1}\right\rangle \otimes\left|j_{2}\left(j_{2}-1\right)\right\rangle+B\left|j_{1}\left(j_{1}-1\right)\right\rangle \otimes\left|j_{2} j_{2}\right\rangle \tag{11.67}
\end{equation*}
$$

Observe that these are the same kets that we ended up with by application of the lowering operator on the topmost element of the first column in our table. Since $\left|j_{1}+j_{2}, j_{1}+j_{2}-1\right\rangle$ and $\left|j_{1}+j_{2}-1, j_{1}+j_{2}-1\right\rangle$ are orthogonal, we can construct our ket for the top of the second column by just seeking such an orthonormal superposition. Consider for example

$$
\begin{align*}
0 & =(a\langle b|+c\langle d|)(A|b\rangle+C|d\rangle) \\
& =a A+c C \tag{11.68}
\end{align*}
$$

With $A=1$ we find that $C=-a / c$, so we have

$$
\begin{align*}
A|b\rangle+C|d\rangle & =|b\rangle-\frac{a}{c}|d\rangle  \tag{11.69}\\
& \sim c|b\rangle-a|d\rangle .
\end{align*}
$$

So we find, for real $a$ and $c$ that

$$
\begin{equation*}
0=(a\langle b|+c\langle d|)(c|b\rangle-a|d\rangle), \tag{11.70}
\end{equation*}
$$

for any orthonormal pair of kets $|a\rangle$ and $|d\rangle$. Using this we find

$$
\begin{align*}
\left|j_{1}+j_{2}-1, j_{1}+j_{2}-1\right\rangle= & \left(\frac{j_{2}}{j_{1}+j_{2}}\right)^{1 / 2}\left|j_{1} j_{1}\right\rangle \otimes\left|j_{2}\left(j_{2}-1\right)\right\rangle \\
& -\left(\frac{j_{1}}{j_{1}+j_{2}}\right)^{1 / 2}\left|j_{1}\left(j_{1}-1\right)\right\rangle \otimes\left|j_{2} j_{2}\right\rangle \tag{11.71}
\end{align*}
$$

This will work, although we could also multiply by any phase factor if desired. Such a choice of phase factors is essentially just a convention.

The Clebsch-Gordon convention. This is the convention we will use, where we

- choose the coefficients to be real.
- require the coefficient of the $m_{1}=j_{1}$ term to be $\geq 0$

This gives us the first state in the second column, and we can proceed to iterate using the lowering operators to get all those values. Moving on to the third column

$$
\begin{equation*}
\left|j_{1}+j_{2}-2, j_{1}+j_{2}-2\right\rangle, \tag{11.72}
\end{equation*}
$$

can only be made up of $\left|j_{1} m_{1}\right\rangle \otimes\left|j_{2} m_{2}\right\rangle$ with $m_{1}+m_{2}=j_{1}+j_{2}-2$. There are now three possibilities

$$
\begin{array}{ll}
m_{1}=j_{1} & m_{2}=j_{2}-2 \\
m_{1}=j_{1}-2 & m_{2}=j_{2}  \tag{11.73}\\
m_{1}=j_{1}-1 & m_{2}=j_{2}-1,
\end{array}
$$

and 2 orthogonality conditions, plus conventions. This is enough to determine the ket in the third column. We can formally write

$$
\begin{equation*}
\left|j m ; j_{1} j_{2}\right\rangle=\sum_{m_{1}, m_{2}}\left|j_{1} m_{1}, j_{2} m_{2}\right\rangle\left\langle j_{1} m_{1}, j_{2} m_{2} \mid j m ; j_{1} j_{2}\right\rangle, \tag{11.74}
\end{equation*}
$$

where

$$
\begin{equation*}
\left|j_{1} m_{1}, j_{2} m_{2}\right\rangle=\left|j_{1} m_{1}\right\rangle \otimes\left|j_{2} m_{2}\right\rangle, \tag{11.75}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\langle j_{1} m_{1}, j_{2} m_{2} \mid j m ; j_{1} j_{2}\right\rangle, \tag{11.76}
\end{equation*}
$$

are the Clebsch-Gordon coefficients, sometimes written as

$$
\begin{equation*}
\left\langle j_{1} m_{1}, j_{2} m_{2} \mid j m\right\rangle . \tag{11.77}
\end{equation*}
$$

## Properties

1. $\left\langle j_{1} m_{1}, j_{2} m_{2} \mid j m\right\rangle \neq 0$ only if $j_{1}-j_{2} \leq j \leq j_{1}+j+2$. This is sometimes called the triangle inequality, depicted in fig. 11.1.


Figure 11.1: Angular momentum triangle inequality.
2. $\left\langle j_{1} m_{1}, j_{2} m_{2} \mid j m\right\rangle \neq 0$ only if $m=m_{1}+m_{2}$.
3. Real (convention).
4. $\left\langle j_{1} j_{1}, j_{2}\left(j-j_{1}\right) \mid j j\right\rangle$ positive (convention again).
5. Proved in the text. If follows that

$$
\begin{equation*}
\left\langle j_{1} m_{1}, j_{2} m_{2} \mid j m\right\rangle=(-1)^{j_{1}+j_{2}-j}\left\langle j_{1}\left(-m_{1}\right), j_{2}\left(-m_{2}\right) \mid j(-m)\right\rangle . \tag{11.78}
\end{equation*}
$$

Note that the $\left\langle j_{1} m_{1}, j_{2} m_{2} \mid j m\right\rangle$ are all real. So, they can be assembled into an orthogonal matrix. Example

$$
\left[\begin{array}{l}
|11\rangle  \tag{11.79}\\
|10\rangle \\
|11\rangle \\
|00\rangle
\end{array}\right]=\left[\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & \frac{1}{\sqrt{2}} & \frac{1}{\sqrt{2}} & 0 \\
0 & 0 & 0 & 1 \\
0 & \frac{1}{\sqrt{2}} & \frac{-1}{\sqrt{2}} & 0
\end{array}\right]\left[\begin{array}{l}
|++\rangle \\
|++\rangle \\
|-+\rangle \\
|--\rangle
\end{array}\right] .
$$

Example. Electrons. Consider the special case of an electron, a spin $1 / 2$ particle with $s=1 / 2$ and $m_{s}= \pm 1 / 2$ where we have

$$
\begin{equation*}
\mathbf{J}=\mathbf{L}+\mathbf{S} \tag{11.80}
\end{equation*}
$$

and

$$
\begin{equation*}
|l m\rangle \otimes\left|\frac{1}{2} m_{s}\right\rangle \tag{11.81}
\end{equation*}
$$

possible values of $j$ are $l \pm 1 / 2$

$$
\begin{equation*}
l \otimes \frac{1}{2}=\left(l+\frac{1}{2}\right) \oplus\left(l-\frac{1}{2}\right) . \tag{11.82}
\end{equation*}
$$

Our table representation is then
Table 11.2: Total angular momentum states for electron.

| $j=l+\frac{1}{2}$ | $j=l-\frac{1}{2}$ |
| :---: | :---: |
| $\left\|l+\frac{1}{2}, l+\frac{1}{2}\right\rangle$ |  |
| $\left\|l+\frac{1}{2}, l+\frac{1}{2}-1\right\rangle$ | $\left\|l-\frac{1}{2}, l-\frac{1}{2}\right\rangle$ |
|  | $\left\lvert\, l-\frac{1}{2}\right.,-\left(l-\frac{1}{2}\right\rangle$ |
| $\left\|l+\frac{1}{2},-\left(l+\frac{1}{2}\right)\right\rangle$ |  |

Here $\left|l+\frac{1}{2}, m\right\rangle$ can only have contributions from

$$
\begin{align*}
& \left|l, m-\frac{1}{2}\right\rangle \otimes\left|\frac{1}{2} \frac{1}{2}\right\rangle \\
& \left|l, m+\frac{1}{2}\right\rangle \otimes\left|\frac{1}{2} \frac{1}{2}\right\rangle \tag{11.83}
\end{align*}
$$

$\left|l-\frac{1}{2}, m\right\rangle$ from the same two. So using this and conventions we can work out (in §28 page 524, of our text [4]).

$$
\begin{align*}
\left|l \pm \frac{1}{2}, m\right\rangle & = \pm \frac{1}{\sqrt{2 l+1}}\left(l+\frac{1}{2} \pm m\right)^{1 / 2}\left|l, m-\frac{1}{2}\right\rangle \times\left|\frac{1}{2} \frac{1}{2}\right\rangle \\
& \pm \frac{1}{\sqrt{2 l+1}}\left(l+\frac{1}{2} \mp m\right)^{1 / 2}\left|l, m+\frac{1}{2}\right\rangle \times\left|\frac{1}{2} \frac{1}{2}\right\rangle \tag{11.84}
\end{align*}
$$

### 11.4 TENSOR OPERATORS.

$\S 29$ of the text.
Recall how we characterized a rotation

$$
\begin{equation*}
\mathbf{r} \rightarrow \mathcal{R}(\mathbf{r}) . \tag{11.85}
\end{equation*}
$$

Here we are using an active rotation as depicted in fig. 11.2. Suppose that


Figure 11.2: active rotation.

$$
\begin{equation*}
[\mathcal{R}(\mathbf{r})]_{i}=\sum_{j} M_{i j} r_{j} \tag{11.86}
\end{equation*}
$$

so that

$$
\begin{equation*}
U=e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{J} / \hbar} \tag{11.87}
\end{equation*}
$$

rotates in the same way. Rotating a ket as in fig. 11.3. Rotating a ket


Figure 11.3: Rotating a wavefunction.
using the prescription

$$
\begin{equation*}
\left|\psi^{\prime}\right\rangle=e^{-i \theta \hat{\mathbf{n}} \cdot \mathbf{J} / \hbar}|\psi\rangle \tag{11.89}
\end{equation*}
$$

and write

$$
\begin{equation*}
\left|\psi^{\prime}\right\rangle=U[M]|\psi\rangle \tag{11.90}
\end{equation*}
$$

Now look at

$$
\begin{equation*}
\langle\psi| O|\psi\rangle, \tag{11.91}
\end{equation*}
$$

and compare with
(*)

$$
\begin{equation*}
\left\langle\psi^{\prime}\right| O\left|\psi^{\prime}\right\rangle=\langle\psi| \frac{1}{U^{\dagger}[M] O U[M]}|\psi\rangle \tag{11.92}
\end{equation*}
$$

We will be looking in more detail at (*).

### 12.1 SETUP.

READING: §28 [4].
Rotating with $U[M]$ as in fig. 12.1 to the new position


Figure 12.1: Rotating a state centered at $F$.

$$
\begin{equation*}
\tilde{r}_{i}=\sum_{j} M_{i j} \bar{r}_{j} \tag{12.1}
\end{equation*}
$$

where

$$
\begin{equation*}
\langle\psi| R_{i}|\psi\rangle=\bar{r}_{i} \tag{12.2}
\end{equation*}
$$

$$
\begin{align*}
\langle\psi| U^{\dagger}[M] R_{i} U[M]|\psi\rangle & =\tilde{r}_{i}=\sum_{j} M_{i j} \bar{r}_{j}  \tag{12.3}\\
& =\langle\psi|\left(U^{\dagger}[M] R_{i} U[M]\right)|\psi\rangle .
\end{align*}
$$

So

$$
\begin{equation*}
U^{\dagger}[M] R_{i} U[M]=\sum_{j} M_{i j} R_{j} \tag{12.4}
\end{equation*}
$$

Any three operators $V_{x}, V_{y}, V_{z}$ that transform according to

$$
\begin{equation*}
U^{\dagger}[M] V_{i} U[M]=\sum_{j} M_{i j} V_{j}, \tag{12.5}
\end{equation*}
$$

form the components of a vector operator.

## 12.2 infinitesimal rotations.

Consider infinitesimal rotations, where we can show (problem set 11, problem 1) that

$$
\begin{equation*}
\left[V_{i}, J_{j}\right]=i \hbar \sum_{k} \epsilon_{i j k} V_{k} . \tag{12.6}
\end{equation*}
$$

Note that for $V_{i}=J_{i}$ we recover the familiar commutator rules for angular momentum, but this also holds for operators $\mathbf{R}, \mathbf{P}, \mathbf{J}, \ldots$ Note that

$$
\begin{equation*}
U^{\dagger}[M]=U\left[M^{-1}\right]=U\left[M^{\mathrm{T}}\right], \tag{12.7}
\end{equation*}
$$

so

$$
\begin{equation*}
U^{\dagger}[M] V_{i} U^{\dagger}[M]=U^{\dagger}\left[M^{\dagger}\right] V_{i} U\left[M^{\dagger}\right]=\sum_{j} M_{j i} V_{j}, \tag{12.8}
\end{equation*}
$$

SO

$$
\begin{equation*}
\langle\psi| V_{i}|\psi\rangle=\langle\psi| U^{\dagger}[M]\left(U[M] V_{i} U^{\dagger}[M]\right) U[M] .|\psi\rangle \tag{12.9}
\end{equation*}
$$

In the same way, suppose we have nine operators

$$
\begin{equation*}
\tau_{i j}, \quad i, j=x, y, z . \tag{12.10}
\end{equation*}
$$

These transform according to

$$
\begin{equation*}
U[M] \tau_{i j} U^{\dagger}[M]=\sum_{l m} M_{l i} M_{m j} \tau_{l m}, \tag{12.11}
\end{equation*}
$$

and we will call these the components of (Cartesian) a second rank tensor operator. Suppose that we have an operator $S$ that transforms as

$$
\begin{equation*}
U[M] S U^{\dagger}[M]=S . \tag{12.12}
\end{equation*}
$$

We will call $S$ a scalar operator.

### 12.3 A PROBLEM.

This all looks good, but it is really not satisfactory. There is a problem.
Suppose that we have a Cartesian tensor operator like this, lets look at the quantity

$$
\begin{align*}
\sum_{i} \tau_{i i} & =\sum_{i} U[M] \tau_{i i} U^{\dagger}[M] \\
& =\sum_{i} \sum_{l m} M_{l i} M_{m i} \tau_{l m} \\
& =\sum_{i} \sum_{l m} M_{l i} M_{i m}^{\mathrm{T}} \tau_{l m}  \tag{12.13}\\
& =\sum_{l m} \delta_{l m} \tau_{l m} \\
& =\sum_{l} \tau_{l l}
\end{align*}
$$

We see buried inside these Cartesian tensors of higher rank there is some simplicity embedded (in this case trace invariance). Who knows what other relationships are also there? We want to work with and extract the buried simplicities, and we will find that the Cartesian way of expressing these tensors is horribly inefficient. What is a representation that does not have any excess information, and is in some sense minimal?

### 12.4 HOW DO WE EXTRACT THESE BURIED SIMPLICITIES?

Recall

$$
\begin{equation*}
U[M]\left|j m^{\prime \prime}\right\rangle \tag{12.14}
\end{equation*}
$$

gives a linear combination of the $\left|j m^{\prime}\right\rangle$.

$$
\begin{align*}
U[M]\left|j m^{\prime \prime}\right\rangle & =\sum_{m^{\prime}}\left|j m^{\prime}\right\rangle\left\langle j m^{\prime}\right| U[M]\left|j m^{\prime \prime}\right\rangle \\
& =\sum_{m^{\prime}}\left|j m^{\prime}\right\rangle D_{m^{\prime} m^{\prime \prime}}^{(j)}[M] . \tag{12.15}
\end{align*}
$$

We have talked about before how these $D_{m^{\prime} m^{\prime \prime}}^{(j)}[M]$ form a representation of the rotation group. These are in fact (not proved here) an irreducible representation.

Look at each element of $D_{m^{\prime} m^{\prime \prime}}^{(j)}[M]$. These are matrices and will be different according to which rotation $M$ is chosen. There is some $M$ for which this element is nonzero. There is no element in this matrix element that is zero for all possible $M$. There are more formal ways to think about this in a group theory context, but this is a physical way to think about this. Think of these as the basis vectors for some eigenket of $J^{2}$.

$$
\begin{align*}
|\psi\rangle & =\sum_{m^{\prime \prime}}\left|j m^{\prime \prime}\right\rangle\left\langle j m^{\prime \prime} \mid \psi\right\rangle  \tag{12.16}\\
& =\sum_{m^{\prime \prime}} \bar{a}_{m^{\prime \prime}}\left|j m^{\prime \prime}\right\rangle
\end{align*}
$$

where

$$
\begin{equation*}
\bar{a}_{m^{\prime \prime}}=\left\langle j m^{\prime \prime} \mid \psi\right\rangle \tag{12.17}
\end{equation*}
$$

So

$$
\begin{align*}
U[M]|\psi\rangle= & =\sum_{m^{\prime}} U[M]\left|j m^{\prime}\right\rangle\left\langle j m^{\prime} \mid \psi\right\rangle \\
& =\sum_{m^{\prime}} U[M]\left|j m^{\prime}\right\rangle \bar{a}_{m^{\prime}} \\
& =\sum_{m^{\prime}, m^{\prime \prime}}\left|j m^{\prime \prime}\right\rangle\left\langle j m^{\prime \prime}\right| U[M]\left|j m^{\prime}\right\rangle \bar{a}_{m^{\prime}}  \tag{12.18}\\
& =\sum_{m^{\prime}, m^{\prime \prime}}\left|j m^{\prime \prime}\right\rangle D_{m^{\prime \prime}, m^{\prime}}^{(j)} \bar{a}_{m^{\prime}} \\
& =\sum_{m^{\prime \prime}} \tilde{a}_{m^{\prime \prime}}\left|j m^{\prime \prime}\right\rangle
\end{align*}
$$

where

$$
\begin{equation*}
\tilde{a}_{m^{\prime \prime}}=\sum_{m^{\prime}} D_{m^{\prime \prime}, m^{\prime}}^{(j)} \bar{a}_{m^{\prime}} \tag{12.19}
\end{equation*}
$$

Recall that

$$
\begin{equation*}
\tilde{r}_{j}=\sum_{j} M_{i j} \bar{r}_{j} \tag{12.20}
\end{equation*}
$$

Define $(2 k+1)$ operators $T_{k} q, q=k, k-1, \cdots-k$ as the elements of a spherical tensor of rank $k$ if

$$
\begin{equation*}
U[M] T_{k}^{q} U^{\dagger}[M]=\sum_{q^{\prime}} D_{q^{\prime} q}^{(j)} T_{k}^{q^{\prime}} \tag{12.21}
\end{equation*}
$$

Here we are looking for a better way to organize things, and it will turn out (not to be proved) that this will be an irreducible way to represent things.

### 12.5 MOTIVATING SPHERICAL TENSORS.

We want to work though some examples of spherical tensors, and how they relate to Cartesian tensors. To do this, a motivating story needs to be told. Let us suppose that $|\psi\rangle$ is a ket for a single particle. Perhaps we are talking about an electron without spin, and write

$$
\begin{align*}
\langle\mathbf{r} \mid \psi\rangle & =Y_{l m}(\theta, \phi) f(r) \\
& =\sum_{m^{\prime \prime}} \bar{a}_{m^{\prime \prime}} Y_{l m^{\prime \prime}}(\theta, \phi), \tag{12.22}
\end{align*}
$$

for $\bar{a}_{m^{\prime \prime}}=\delta_{m^{\prime \prime} m}$ and after dropping $f(r)$. So

$$
\begin{equation*}
\langle\mathbf{r}| U[M]|\psi\rangle=\sum_{m^{\prime \prime}} \sum_{m^{\prime}} D_{m^{\prime \prime} m}^{(j)} \bar{a}_{m^{\prime}} Y_{l m^{\prime \prime}}(\theta, \phi) \tag{12.23}
\end{equation*}
$$

We are writing this in this particular way to make a point. Now also assume that

$$
\begin{equation*}
\langle\mathbf{r} \mid \psi\rangle=Y_{l m}(\theta, \phi), \tag{12.24}
\end{equation*}
$$

so we find

$$
\begin{align*}
&\langle\mathbf{r}| U[M]|\psi\rangle=\sum_{m^{\prime \prime}} Y_{l m^{\prime \prime}}(\theta, \phi) D_{m^{\prime \prime} m}^{(j)}  \tag{12.25}\\
&=Y_{l m}(\theta, \phi), \\
& Y_{l m}(\theta, \phi)=Y_{l m}(x, y, z) \tag{12.26}
\end{align*}
$$

so

$$
\begin{equation*}
Y_{l m}^{\prime}(x, y, z)=\sum_{m^{\prime \prime}} Y_{l m^{\prime \prime}}(x, y, z) D_{m^{\prime \prime} m}^{(j)} \tag{12.27}
\end{equation*}
$$

Now consider the spherical harmonic as an operator $Y_{\operatorname{lm}}(X, Y, Z)$

$$
\begin{equation*}
U[M] Y_{l m}(X, Y, Z) U^{\dagger}[M]=\sum_{m^{\prime \prime}} Y_{l m^{\prime \prime}}(X, Y, Z) D_{m^{\prime \prime} m}^{(j)} \tag{12.28}
\end{equation*}
$$

So this is a way to generate spherical tensor operators of rank $0,1,2, \cdots$.

### 12.6 SPHERICAL TENSORS (CONT.)

READING: §29 of [4].

Definition. Any $(2 k+1)$ operator $T(k, q), q=-k, \cdots, k$ are the elements of a spherical tensor of rank $k$ if

$$
\begin{equation*}
U[M] T(k, q) U^{-1}[M]=\sum_{q^{\prime}} T\left(k, q^{\prime}\right) D_{q q^{\prime}}^{(k)} \tag{12.29}
\end{equation*}
$$

where $D_{q q^{\prime}}^{(k)}$ was the matrix element of the rotation operator

$$
\begin{equation*}
D_{q q^{\prime}}^{(k)}=\left\langle k q^{\prime}\right| U[M]\left|k q^{\prime \prime}\right\rangle \tag{12.30}
\end{equation*}
$$

So, if we have a Cartesian vector operator with components $V_{x}, V_{y}, V_{z}$ then we can construct a corresponding spherical vector operator

$$
\begin{array}{lll}
T(1,1) & =-\frac{V_{x}+i V_{y}}{\sqrt{2}} & \equiv V_{+1} \\
T(1,0) & =V_{z} & \equiv V_{0}  \tag{12.31}\\
T(1,-1) & =-\frac{V_{x}-i V_{y}}{\sqrt{2}} & \equiv V_{-1}
\end{array}
$$

By considering infinitesimal rotations we can come up with the commutation relations between the angular momentum operators

$$
\begin{align*}
{\left[J_{ \pm}, T(k, q)\right] } & =\hbar \sqrt{(k \mp q)(k \pm q+1)} T(k, q \pm 1)  \tag{12.32}\\
{\left[J_{z}, T(k, q)\right] } & =\hbar q T(k, q)
\end{align*}
$$

Note that the text in (29.15) defines these, whereas in class these were considered consequences of eq. (12.29), once infinitesimal rotations were used. Recall that these match our angular momentum raising and lowering identities

$$
\begin{align*}
J_{ \pm}|k q\rangle & =\hbar \sqrt{(k \mp q)(k \pm q+1)}|k, q \pm 1\rangle  \tag{12.33}\\
J_{z}|k q\rangle & =\hbar q|k, q\rangle .
\end{align*}
$$

Consider two problems

$$
\begin{array}{lll}
T(k, q) & & |k q\rangle \\
{\left[J_{ \pm}, T(k, q)\right]} & \leftrightarrow & J_{ \pm}|k q\rangle  \tag{12.34}\\
{\left[J_{z}, T(k, q)\right]} & \leftrightarrow & J_{z}|k q\rangle .
\end{array}
$$

We have a correspondence between the spherical tensors and angular momentum kets

| $T_{1}\left(k_{1}, q_{1}\right)$ | $q_{1}=-k_{1}, \cdots, k_{1}$ | $\left\|k_{1} q_{1}\right\rangle$ | $\left\|k_{2} q_{2}\right\rangle$ |
| :--- | :--- | :--- | :--- |
| $T_{2}\left(k_{2}, q_{2}\right)$ | $q_{2}=-k_{2}, \cdots, k_{2}$ | $q_{1}=-k_{1}, \cdots k_{1}$ | $q_{2}=-k_{2}, \cdots k_{2}$. |

So, as we can write for angular momentum

These are the C.G coefficients

$$
\begin{align*}
|k q\rangle & =\sum_{q_{1}, q_{2}}\left|k_{1}, q_{1}\right\rangle\left|k_{2}, q_{2}\right\rangle\left\langle k_{1} q_{1} k_{2} q_{2} \mid k q\right\rangle  \tag{12.36}\\
\left|k_{1} q_{1} ; k_{2} q_{2}\right\rangle & =\sum_{k, q^{\prime}}\left|k q^{\prime}\right\rangle\left\langle k q^{\prime} \mid k_{1} q_{1} k_{2} q_{2}\right\rangle .
\end{align*}
$$

We also have for spherical tensors

$$
\begin{align*}
T(k, q) & =\sum_{q_{1}, q_{2}} T_{1}\left(k_{1}, q_{1}\right) T_{2}\left(k_{2}, q_{2}\right)\left\langle k_{1} q_{1} k_{2} q_{2} \mid k q\right\rangle \\
T_{1}\left(k_{1}, q_{1}\right) T_{2}\left(k_{2}, q_{2}\right) & =\sum_{k, q^{\prime}} T\left(k, q^{\prime}\right)\left\langle k q^{\prime} \mid k_{1} q_{1} k_{2} q_{2}\right\rangle \tag{12.37}
\end{align*}
$$

Can form eigenstates $|k q\rangle$ of (total angular momentum) ${ }^{2}$ and (z-comp of the total angular momentum). FIXME: this will not be proven, but we are strongly suggested to try this ourselves.

$$
\begin{array}{lll}
\text { spherical tensor (3) } & \leftrightarrow & \text { Cartesian vector (3) }  \tag{12.38}\\
\text { (spherical vector)(spherical vector) } & & \text { Cartesian tensor. }
\end{array}
$$

We can check the dimensions for a spherical tensor decomposition into rank 0 , rank 1 and rank 2 tensors.

| spherical tensor rank 0 | (1) | (Cartesian vector)(Cartesian vector) |
| :--- | :--- | :--- |
| spherical tensor rank 1 | (3) | $(3)(3)$ |
| spherical tensor rank 2 | (5) | 9 |
| dimension check sum | 9 |  |.

Or in the direct product and sum shorthand

$$
\begin{equation*}
1 \otimes 1=0 \oplus 1 \oplus 2 \tag{12.40}
\end{equation*}
$$

Note that this is just like problem 4 in problem set 10 where we calculated the CG kets for the $1 \otimes 1=0 \oplus 1 \oplus 2$ decomposition starting from kets $|1 m\rangle\left|1 m^{\prime}\right\rangle$.
|21> |11>
$|20\rangle|10\rangle|00\rangle$
$|2 \overline{1}\rangle \quad|1 \overline{1}\rangle$
$|2 \overline{2}\rangle$
Example. How about a Cartesian tensor of rank 3?

$$
\begin{equation*}
A_{i j k} \tag{12.42}
\end{equation*}
$$

$$
\begin{align*}
& 1 \otimes 1 \otimes 1=1 \otimes(0 \oplus 1 \oplus 2) \\
& =(1 \otimes 0) \oplus(1 \otimes 1) \oplus(1 \otimes 2) \\
& =\begin{array}{l}
1 \oplus(0 \oplus 1 \oplus 2) \oplus(3 \oplus 2 \oplus 1) \\
3+1+3+5+7+5+3=27 .
\end{array} \tag{12.43}
\end{align*}
$$

Why bother? Consider a tensor operator $T(k, q)$ and an eigenket of angular momentum $|\alpha j m\rangle$, where $\alpha$ is a degeneracy index. Look at

$$
\begin{align*}
T(k, q)|\alpha j m\rangle U[M] T(k, q)|\alpha j m\rangle & =U[M] T(k, q) U^{\dagger}[M] U[M]|\alpha j m\rangle \\
& =\sum_{q^{\prime} m^{\prime}} D_{q q^{\prime}}^{(k)} D_{m m^{\prime}}^{(j)} T\left(k, q^{\prime}\right)\left|\alpha j m^{\prime}\right\rangle . \tag{12.44}
\end{align*}
$$

This transforms like $|k q\rangle \otimes|j m\rangle$. We can say immediately

$$
\begin{equation*}
\left\langle\alpha^{\prime} j^{\prime} m^{\prime}\right| T(k, q)|\alpha j m\rangle=0, \tag{12.45}
\end{equation*}
$$

unless

$$
\begin{align*}
|k-j| & \leq j^{\prime} \leq k+j  \tag{12.46}\\
m^{\prime} & =m+q .
\end{align*}
$$

This is the "selection rule".
Examples.

- Scalar $T(0,0)$

$$
\begin{equation*}
\left\langle\alpha^{\prime} j^{\prime} m^{\prime}\right| T(0,0)|\alpha j m\rangle=0, \tag{12.47}
\end{equation*}
$$

unless $j=j^{\prime}$ and $m=m^{\prime}$.

- $V_{x}, V_{y}, V_{z}$. What are the non-vanishing matrix elements?

$$
\begin{align*}
& V_{x}=\frac{V_{-1}-V_{+1}}{\sqrt{2}}, \cdots  \tag{12.48}\\
& \left\langle\alpha^{\prime} j^{\prime} m^{\prime}\right| V_{x, y}|\alpha j m\rangle=0, \tag{12.49}
\end{align*}
$$

unless

$$
\begin{align*}
|j-1| & \leq j^{\prime} \leq j+1 \\
m^{\prime} & =m \pm 1 . \tag{12.50}
\end{align*}
$$

$$
\begin{equation*}
\left\langle\alpha^{\prime} j^{\prime} m^{\prime}\right| V_{z}|\alpha j m\rangle=0, \tag{12.51}
\end{equation*}
$$

unless

$$
\begin{align*}
|j-1| & \leq j^{\prime} \leq j+1 \\
m^{\prime} & =m . \tag{12.52}
\end{align*}
$$

Very generally one can prove (the Wigner-Eckart theory in the text §29.3)

$$
\begin{equation*}
\left\langle\alpha_{2} j_{2} m_{2}\right| T(k, q)\left|\alpha_{1} j_{1} m_{1}\right\rangle=\left\langle\alpha_{2} j_{2}\right| T(k)\left|\alpha_{1} j_{1}\right\rangle \cdot\left\langle j_{2} m_{2} \mid k q_{1} ; j_{1} m_{1}\right\rangle, \tag{12.53}
\end{equation*}
$$

where we split into a "reduced matrix element" describing the "physics", and the CG coefficient for "geometry" respectively.

### 13.1 SETUP.

READING: $\S 19, \S 20$ of the text [4].
Figure 13.1 shows a simple classical picture of a two particle scattering collision We will focus on point particle elastic collisions (no energy lost


Figure 13.1: classical collision of particles.
in the collision). With particles of mass $m_{1}$ and $m_{2}$ let the total mass be

$$
\begin{equation*}
M=m_{1}+m_{2}, \tag{13.1}
\end{equation*}
$$

with a corresponding reduced mass of

$$
\begin{equation*}
\frac{1}{\mu}=\frac{1}{m_{1}}+\frac{1}{m_{2}}, \tag{13.2}
\end{equation*}
$$

so that interaction due to a potential $V\left(\mathbf{r}_{1}-\mathbf{r}_{2}\right)$ that depends on the difference in position $\mathbf{r}=\mathbf{r}_{1}-\mathbf{r}$ has, in the center of mass frame, the Hamiltonian

$$
\begin{equation*}
H=\frac{\mathbf{p}^{2}}{2 \mu}+V(\mathbf{r}) . \tag{13.3}
\end{equation*}
$$

In the classical picture we would investigate the scattering radius $r_{0}$ associated with the impact parameter $\rho$ as depicted in fig. 13.2.


Figure 13.2: Classical scattering radius and impact parameter.

### 13.2 1d QM SCATTERING. NO POTENTIAL WAVE PACKET TIME EVOLUTION.

Now lets move to the QM picture where we assume that we have a particle that can be represented as a wave packet as in fig. 13.3. First without


Figure 13.3: Wave packet for a particle wavefunction $\operatorname{Re}(\psi(x, 0))$.
any potential $V(x)=0$, lets consider the evolution. Our position and momentum space representations are related by

$$
\begin{equation*}
\int|\psi(x, t)|^{2} d x=1=\int|\psi(p, t)|^{2} d p \tag{13.4}
\end{equation*}
$$

and by Fourier transform

$$
\begin{equation*}
\psi(x, t)=\int \frac{d p}{\sqrt{2 \pi \hbar}} \bar{\psi}(p, t) e^{i p x / \hbar} \tag{13.5}
\end{equation*}
$$

Schrödinger's equation takes the form

$$
\begin{equation*}
i \hbar \frac{\partial \psi(x, t)}{\partial t}=-\frac{\hbar^{2}}{2 \mu} \frac{\partial^{2} \psi(x, t)}{\partial x^{2}} \tag{13.6}
\end{equation*}
$$

or more simply in momentum space

$$
\begin{equation*}
i \hbar \frac{\partial \bar{\psi}(p, t)}{\partial t}=\frac{p^{2}}{2 \mu} \frac{\partial^{2} \bar{\psi}(p, t)}{\partial x^{2}} . \tag{13.7}
\end{equation*}
$$

Rearranging to integrate we have

$$
\begin{equation*}
\frac{\partial \bar{\psi}}{\partial t}=-\frac{i p^{2}}{2 \mu \hbar} \bar{\psi}, \tag{13.8}
\end{equation*}
$$

and integrating

$$
\begin{equation*}
\ln \bar{\psi}=-\frac{i p^{2} t}{2 \mu \hbar}+\ln C \tag{13.9}
\end{equation*}
$$

or

$$
\begin{equation*}
\bar{\psi}=C e^{-\frac{i p^{2} t}{2 \mu \hbar}}=\bar{\psi}(p, 0) e^{-\frac{i p^{2} t}{2 \mu \hbar}} . \tag{13.10}
\end{equation*}
$$

Time evolution in momentum space for the free particle changes only the phase of the wavefunction, the momentum probability density of that particle.

Fourier transforming, we find our position space wavefunction to be

$$
\begin{equation*}
\psi(x, t)=\int \frac{d p}{\sqrt{2 \pi \hbar}} \bar{\psi}(p, 0) e^{i p x / \hbar} e^{-i p^{2} t / 2 \mu \hbar} . \tag{13.11}
\end{equation*}
$$

To clean things up, write

$$
\begin{equation*}
p=\hbar k \tag{13.12}
\end{equation*}
$$

for

$$
\begin{equation*}
\left.\psi(x, t)=\int \frac{d k}{\sqrt{2 \pi}} a(k, 0)\right) e^{i k x} e^{-i \hbar k^{2} t / 2 \mu} \tag{13.13}
\end{equation*}
$$

where

$$
\begin{equation*}
a(k, 0)=\sqrt{\hbar} \bar{\psi}(p, 0) . \tag{13.14}
\end{equation*}
$$

Putting

$$
\begin{equation*}
a(k, t)=a(k, 0) e^{-i \hbar k^{2} / 2 \mu} \tag{13.15}
\end{equation*}
$$

we have

$$
\begin{equation*}
\left.\psi(x, t)=\int \frac{d k}{\sqrt{2 \pi}} a(k, t)\right) e^{i k x} \tag{13.16}
\end{equation*}
$$

Observe that we have

$$
\begin{equation*}
\int d k|a(k, t)|^{2}=\int d p|\bar{\psi}(p, t)|^{2}=1 \tag{13.17}
\end{equation*}
$$

## 13.3 a gaussian wave packet.

Suppose that we have, as depicted in fig. 4.8, a Gaussian wave packet of the form

$$
\begin{equation*}
\psi(x, 0)=\left(\pi \Delta^{2}\right)^{1 / 4} e^{i k_{0} x} e^{-x^{2} / 2 \Delta^{2}} \tag{13.18}
\end{equation*}
$$

This is actually a minimum uncertainty packet with

$$
\begin{align*}
& \Delta x=\frac{\Delta}{\sqrt{2}} \\
& \Delta p=\frac{\hbar}{\Delta \sqrt{2}} . \tag{13.19}
\end{align*}
$$

Taking Fourier transforms we have

$$
\begin{align*}
& a(k, 0)=\left(\frac{\Delta^{2}}{\pi}\right)^{1 / 4} e^{-\left(k-k_{0}\right)^{2} \Delta^{2} / 2},  \tag{13.20}\\
& a(k, t)=\left(\frac{\Delta^{2}}{\pi}\right)^{1 / 4} e^{-\left(k-k_{0}\right)^{2} \Delta^{2} / 2} e^{-i \hbar k^{2} t / 2 \mu} \equiv \alpha(k, t)
\end{align*}
$$

For $t>0$ our wave packet will start moving and spreading as in fig. 13.4.


Figure 13.4: moving spreading Gaussian packet.

## 13.4 with a potential.

Now "switch on" a potential, still assuming a wave packet representation for the particle. With a positive (repulsive) potential as in fig. 13.5, at a time long before the interaction of the wave packet with the potential we can visualize the packet as heading towards the barrier. After some time


Figure 13.5: QM wave packet prior to interaction with repulsive potential.
long after the interaction, classically for this sort of potential where the particle kinetic energy is less than the barrier "height", we would have total reflection. In the QM case, we have seen before that we will have a reflected and a transmitted portion of the wave packet as depicted in fig. 13.6. Even if the particle kinetic energy is greater than the barrier


Figure 13.6: QM wave packet long after interaction with repulsive potential.
height, or for a negative potential (both such potentials are depicted in fig. 13.7), then we can still have a reflected component. Consider the


Figure 13.7: Kinetic energy greater than two potential energy distributions.
probability for the particle to be found anywhere long after the interac-
tion, summing over the transmitted and reflected wave functions, we have

$$
\begin{align*}
1 & =\int\left|\psi_{r}+\psi_{t}\right|^{2}  \tag{13.21}\\
& =\int\left|\psi_{r}\right|^{2}+\int\left|\psi_{t}\right|^{2}+2 \mathfrak{R} \int \psi_{r}^{*} \psi_{t}
\end{align*}
$$

Observe that long after the interaction the cross terms in the probabilities will vanish because they are non-overlapping, leaving just the probably densities for the transmitted and reflected probably densities independently.

We define

$$
\begin{align*}
& T=\int\left|\psi_{t}(x, t)\right|^{2} d x,  \tag{13.22}\\
& R=\int\left|\psi_{r}(x, t)\right|^{2} d x .
\end{align*}
$$

The objective of most of our scattering problems will be the calculation of these probabilities and the comparisons of their ratios.

Question. Can we have more than one wave packet reflect off. Yes, we could have multiple wave packets for both the reflected and the transmitted portions. For example, if the potential has some internal structure there could be internal reflections before anything emerges on either side and things could get quite messy.

## 13.5 considering the time independent case temporarily.

We are going to work through something that is going to seem at first to be completely unrelated. We will (eventually) see that this can be applied to this problem, so a bit of patience will be required.

We will be using the time independent Schrödinger equation

$$
\begin{equation*}
-\frac{\hbar^{2}}{2 \mu} \psi_{k}^{\prime \prime}(x)=V(x) \psi_{k}(x)=E \psi_{k}(x) \tag{13.23}
\end{equation*}
$$

where we have added a subscript $k$ to our wave function with the intention (later) of allowing this to vary. For "future use" we define for $k>0$

$$
\begin{equation*}
E=\frac{\hbar^{2} k^{2}}{2 \mu} . \tag{13.24}
\end{equation*}
$$

Consider a potential as in fig. 13.8, where $V(x)=0$ for $x>x_{2}$ and $x<x_{1}$. We will not have bound states here (repulsive potential). There will be


Figure 13.8: potential zero outside of a specific region.
many possible solutions, but we want to look for a solution that is of the form

$$
\begin{equation*}
\psi_{k}(x)=C e^{i k x}, \quad x>x_{2} \tag{13.25}
\end{equation*}
$$

Suppose $x=x_{3}>x_{2}$, we have

$$
\begin{align*}
& \psi_{k}\left(x_{3}\right)=C e^{i k x_{3}} \\
& \left.\frac{d \psi_{k}}{d x}\right|_{x=x_{3}}=i k C e^{i k x_{3}} \equiv \phi_{k}\left(x_{3}\right)  \tag{13.27}\\
& \left.\frac{d^{2} \psi_{k}}{d x^{2}}\right|_{x=x_{3}}=-k^{2} C e^{i k x_{3}} . \tag{13.28}
\end{align*}
$$

Defining

$$
\begin{equation*}
\phi_{k}(x)=\frac{d \psi_{k}}{d x}, \tag{13.29}
\end{equation*}
$$

we write Schrödinger's equation as a pair of coupled first order equations

$$
\begin{align*}
\frac{d \psi_{k}}{d x} & =\phi_{k}(x) \\
-\frac{\hbar^{2}}{2 \mu} \frac{d \phi_{k}(x)}{d x} & =-V(x) \psi_{k}(x)+\frac{\hbar^{2} k^{2}}{2 \mu} \psi_{k}(x) \tag{13.30}
\end{align*}
$$

At this $x=x_{3}$ specifically, we "know" both $\phi_{k}\left(x_{3}\right)$ and $\psi_{k}\left(x_{3}\right)$ and have

$$
\begin{align*}
\left.\frac{d \psi_{k}}{d x}\right|_{x_{3}} & =\phi_{k}(x) \\
-\left.\frac{\hbar^{2}}{2 \mu} \frac{d \phi_{k}(x)}{d x}\right|_{x_{3}} & =-V\left(x_{3}\right) \psi_{k}\left(x_{3}\right)+\frac{\hbar^{2} k^{2}}{2 \mu} \psi_{k}\left(x_{3}\right) . \tag{13.31}
\end{align*}
$$

This allows us to find both

$$
\begin{gather*}
\left.\frac{d \psi_{k}}{d x}\right|_{x_{3}} \\
\left.\frac{d \phi_{k}(x)}{d x}\right|_{x_{3}} \tag{13.32}
\end{gather*}
$$

then proceed to numerically calculate $\phi_{k}(x)$ and $\psi_{k}(x)$ at neighboring points $x=x_{3}+\epsilon$. Essentially, this allows us to numerically integrate backwards from $x_{3}$ to find the wave function at previous points for any sort of potential.

### 13.6 Recap.

READING: $\S 19, \S 20$ of the text [4].
We used a positive potential of the form of fig. 13.9.


Figure 13.9: A bounded positive potential.

$$
\begin{equation*}
-\frac{\hbar^{2}}{2 \mu} \frac{\partial^{2} \psi_{k}(x)}{\partial x^{2}}+V(x) \psi_{k}(x)=\frac{\hbar^{2} k^{2}}{2 \mu} . \tag{13.33}
\end{equation*}
$$

For $x \geq x_{3}$

$$
\begin{equation*}
\psi_{k}(x)=C e^{i k x} \tag{13.34}
\end{equation*}
$$

$$
\begin{equation*}
\phi_{k}(x)=\frac{d \psi_{k}(x)}{d x} \tag{13.35}
\end{equation*}
$$

For $x \geq x_{3}$

$$
\begin{equation*}
\phi_{k}(x)=i k C e^{i k x} \tag{13.36}
\end{equation*}
$$

$$
\begin{align*}
\frac{d \psi_{k}(x)}{d x} & =\phi_{k}(x) \\
-\frac{\hbar^{2}}{2 \mu} \frac{d \phi_{k}(x)}{d x} & =-V(x) \psi_{k}(x)+\frac{\hbar^{2} k^{2}}{2 \mu} \tag{13.37}
\end{align*}
$$

We will integrate these equations back to $x_{1}$.
For $x \leq x_{1}$

$$
\begin{equation*}
\psi_{k}(x)=A e^{i k x}+B e^{-i k x} \tag{13.38}
\end{equation*}
$$

where both $A$ and $B$ are proportional to $C$, dependent on $k$. There are cases where we can solve this analytically (one of these is on our problem set). Alternatively, write as (so long as $A \neq 0$ )

$$
\begin{align*}
\psi_{k}(x) & \rightarrow e^{i k x}+\beta_{k} e^{-i k x} & & \text { for } x<x_{1}  \tag{13.39}\\
& \rightarrow \gamma_{k} e^{i k x} & & \text { for } x>x_{2}
\end{align*}
$$

Now want to consider the problem of no potential in the interval of interest, and our window bounded potential as in fig. 13.10, where we model


Figure 13.10: Wave packet in free space and with positive potential.
our particle as a wave packet as we found can have the Fourier transform description, for $t_{\text {initial }}<0$, of

$$
\begin{equation*}
\psi\left(x, t_{\text {initial }}\right)=\int \frac{d k}{\sqrt{2 \pi}} \alpha\left(k, t_{\text {initial }}\right) e^{i k x} . \tag{13.40}
\end{equation*}
$$

Returning to the same coefficients, the solution of the Schrödinger eqn for problem with the potential eq. (13.39) For $x \leq x_{1}$,

$$
\begin{equation*}
\psi(x, t)=\psi_{i}(x, t)+\psi_{r}(x, t) \tag{13.41}
\end{equation*}
$$

where as illustrated in fig. 13.11, we have


Figure 13.11: Reflection and transmission of wave packet.

$$
\begin{align*}
\psi_{i}(x, t) & =\int \frac{d k}{\sqrt{2 \pi}} \alpha\left(k, t_{\text {initial }}\right) e^{i k x} \\
\psi_{r}(x, t) & =\int \frac{d k}{\sqrt{2 \pi}} \alpha\left(k, t_{\text {initial }}\right) \beta_{k} e^{-i k x} \tag{13.42}
\end{align*}
$$

For $x>x_{2}$

$$
\begin{equation*}
\psi(x, t)=\psi_{t}(x, t) \tag{13.43}
\end{equation*}
$$

and

$$
\begin{equation*}
\psi_{t}(x, t)=\int \frac{d k}{\sqrt{2 \pi}} \alpha\left(k, t_{\text {initial }}\right) \gamma_{k} e^{i k x} \tag{13.44}
\end{equation*}
$$

Look at

$$
\begin{equation*}
\psi_{r}(x, t)=\chi(-x, t) \tag{13.45}
\end{equation*}
$$

where

$$
\begin{align*}
\chi(x, t) & =\int \frac{d k}{\sqrt{2 \pi}} \alpha\left(k, t_{\text {initial }}\right) \beta_{k} e^{i k x} \\
& \approx \beta_{k_{0}} \int \frac{d k}{\sqrt{2 \pi}} \alpha\left(k, t_{\text {initial }}\right) e^{i k x} \tag{13.46}
\end{align*}
$$

for $t=t_{\text {initial }}$, this is nonzero for $x<x_{1}$. So for $x<x_{1}$

$$
\begin{equation*}
\psi_{r}\left(x, t_{\text {initial }}\right)=0 \tag{13.47}
\end{equation*}
$$

In the same way, for $x>x_{2}$

$$
\begin{equation*}
\psi_{t}\left(x, t_{\text {initial }}\right)=0 \tag{13.48}
\end{equation*}
$$

What has not been proved is that the wavefunction is also zero in the [ $x_{1}, x_{2}$ ] interval.

Summarizing. For $t=t_{\text {initial }}$

$$
\psi\left(x, t_{\text {initial }}\right)= \begin{cases}\int \frac{d k}{\sqrt{2 \pi}} \alpha\left(k, t_{\text {initial }}\right) e^{i k x} & \text { for } x<x_{1}  \tag{13.49}\\ 0 & \text { for } x>x_{2}, x>x_{1}\end{cases}
$$

Note that we didn't prove the $x>x_{1}$ case. For $t=t_{\text {final }}$

$$
\psi\left(x, t_{\mathrm{final}}\right) \rightarrow \begin{cases}\int \frac{d k}{\sqrt{2 \pi}} \beta_{k} \alpha\left(k, t_{\mathrm{final}}\right) e^{-i k x} & \text { for } x<x_{1}  \tag{13.50}\\ 0 & x \in\left[x_{1}, x_{2}\right] \\ \int \frac{d k}{\sqrt{2 \pi}} \gamma_{k} \alpha\left(k, t_{\mathrm{final}}\right) e^{i k x} & \text { for } x>x_{2}\end{cases}
$$

The probability of reflection is

$$
\begin{equation*}
\int\left|\psi_{r}\left(x, t_{\mathrm{final}}\right)\right|^{2} d x \tag{13.51}
\end{equation*}
$$

If we have a sufficiently localized packet, we can form a first order approximation around the peak of $\beta_{k}$ (FIXME: or is this a sufficiently localized response to the potential on reflection?)

$$
\begin{equation*}
\psi_{r}\left(x, t_{\mathrm{final}}\right) \approx \beta_{k_{0}} \int \frac{d k}{\sqrt{2 \pi}} \alpha\left(k, t_{\mathrm{final}}\right) e^{-i k x} \tag{13.52}
\end{equation*}
$$

So

$$
\begin{equation*}
\int\left|\psi_{r}\left(x, t_{\mathrm{final}}\right)\right|^{2} d x \approx\left|\beta_{k_{0}}\right|^{2} \equiv R \tag{13.53}
\end{equation*}
$$

The probability of transmission is

$$
\begin{equation*}
\int\left|\psi_{t}\left(x, t_{\mathrm{final}}\right)\right|^{2} d x \tag{13.54}
\end{equation*}
$$

Again, assuming a small spread in $\gamma_{k}$, with $\gamma_{k} \approx \gamma_{k_{0}}$ for some $k_{0}$

$$
\begin{equation*}
\psi_{t}\left(x, t_{\mathrm{final}}\right) \approx \gamma_{k_{0}} \int \frac{d k}{\sqrt{2 \pi}} \alpha\left(k, t_{\mathrm{final}}\right) e^{i k x} \tag{13.55}
\end{equation*}
$$

we have for $x>x_{2}$

$$
\begin{equation*}
\int\left|\psi_{t}\left(x, t_{\mathrm{final}}\right)\right|^{2} d x \approx\left|\gamma_{k_{0}}\right|^{2} \equiv T \tag{13.56}
\end{equation*}
$$

By constructing the wave packets in this fashion we get as a side effect the solution of the scattering problem. The

$$
\begin{align*}
& \psi_{k}(x) \rightarrow e^{i k x}+\beta_{k} e^{-i k x} \\
& \gamma_{k} e^{i k x} \tag{13.57}
\end{align*}
$$

are called asymptotic in states. Their physical applicability is only once we have built wave packets out of them.

### 14.1 SETUP.

READING: $\S 20$, and $\S 4.8$ of our text [4]. For a potential $V(\mathbf{r}) \approx 0$ for $r>r_{0}$ as in fig. 14.1. From 1D we have learned to build up solutions


Figure 14.1: Radially bounded spherical potential.
from time independent solutions (non normalizable). Consider an incident wave

$$
\begin{equation*}
e^{i \mathbf{k} \cdot \mathbf{r}}=e^{i k \hat{\mathbf{n}} \cdot \mathbf{r}} \tag{14.1}
\end{equation*}
$$

This is a solution of the time independent Schrödinger equation

$$
\begin{equation*}
-\frac{\hbar^{2}}{2 \mu} \nabla^{2} e^{i \mathbf{k} \cdot \mathbf{r}}=E e^{i \mathbf{k} \cdot \mathbf{r}} \tag{14.2}
\end{equation*}
$$

where

$$
\begin{equation*}
E=\frac{\hbar^{2} \mathbf{k}^{2}}{2 \mu} \tag{14.3}
\end{equation*}
$$

In the presence of a potential expect scattered waves.
Consider scattering off of a positive potential as depicted in fig. 14.2. Here we have $V(r)=0$ for $r>r_{0}$. The wave function

$$
\begin{equation*}
e^{i k \hat{\mathbf{n}} \cdot \mathbf{r}} \tag{14.4}
\end{equation*}
$$

is found to be a solution of the free particle Schrödinger equation

$$
\begin{equation*}
-\frac{\hbar^{2}}{2 \mu} \nabla^{2} e^{i k \hat{\mathbf{n}} \cdot \mathbf{r}}=\frac{\hbar^{2} \mathbf{k}^{2}}{2 \mu} e^{i k \hat{\mathbf{n}} \cdot \mathbf{r}} \tag{14.5}
\end{equation*}
$$



Figure 14.2: Radially bounded potential.
14.2 SEEKING A POST SCATTERING SOLUTION AWAY FROM THE POTENtial.

What other solutions can be found for $r>r_{0}$, where our potential $V(r)=$ 0 ? We are looking for $\Phi(\mathbf{r})$ such that

$$
\begin{equation*}
-\frac{\hbar^{2}}{2 \mu} \nabla^{2} \Phi(r)=\frac{\hbar^{2} \mathbf{k}^{2}}{2 \mu} \Phi(r) \tag{14.6}
\end{equation*}
$$

What can we find? We split our Laplacian into radial and angular components as we did for the hydrogen atom

$$
\begin{equation*}
-\frac{\hbar^{2}}{2 \mu} \frac{\partial^{2}}{\partial r^{2}}(r \Phi(\mathbf{r}))+\frac{\mathcal{L}^{2}}{2 \mu r^{2}} \Phi(\mathbf{r})=E \Phi(\mathbf{r}) \tag{14.7}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathcal{L}^{2}=-\hbar^{2}\left(\frac{\partial^{2}}{\partial \theta^{2}}+\frac{1}{\tan \theta} \frac{\partial}{\partial \theta}+\frac{1}{\sin ^{2} \theta} \frac{\partial^{2}}{\partial \phi^{2}}\right) \tag{14.8}
\end{equation*}
$$

Assuming a solution of

$$
\begin{equation*}
\Phi(\mathbf{r})=R(r) Y_{l}^{m}(\theta, \phi) \tag{14.9}
\end{equation*}
$$

and noting that

$$
\begin{equation*}
\mathcal{L}^{2} Y_{l}^{m}(\theta, \phi)=\hbar^{2} l(l+1) Y_{l}^{m}(\theta, \phi) \tag{14.10}
\end{equation*}
$$

we find that our radial equation becomes

$$
\begin{equation*}
-\frac{\hbar^{2}}{2 \mu r} \frac{\partial^{2}}{\partial r^{2}}(r R(r))+\frac{\hbar^{2} l(l+1)}{2 \mu r^{2}} R(r)=E R(r)=\frac{\hbar^{2} k^{2}}{2 \mu} R(r) \tag{14.11}
\end{equation*}
$$

Writing

$$
\begin{equation*}
R(r)=\frac{u(r)}{r} \tag{14.12}
\end{equation*}
$$

we have

$$
\begin{equation*}
-\frac{\hbar^{2}}{2 \mu r} \frac{\partial^{2} u(r)}{\partial r^{2}}+\frac{\hbar^{2} l(l+1)}{2 \mu r} u(r)=\frac{\hbar^{2} k^{2}}{2 \mu} \frac{u(r)}{r} \tag{14.13}
\end{equation*}
$$

or

$$
\begin{equation*}
\left(\frac{d^{2}}{d r^{2}}+k^{2}-\frac{l(l+1)}{r^{2}}\right) u(r)=0 \tag{14.14}
\end{equation*}
$$

Writing $\rho=k r$, we have

$$
\begin{equation*}
\left(\frac{d^{2}}{d \rho^{2}}+1-\frac{l(l+1)}{\rho^{2}}\right) u(r)=0 . \tag{14.15}
\end{equation*}
$$

### 14.3 THE RADIAL EQUATION AND ITS SOLUTION.

With a last substitution of $u(r)=U(k r)=U(\rho)$, and introducing an explicit $l$ suffix on our eigenfunction $U(\rho)$ we have

$$
\begin{equation*}
\left(-\frac{d^{2}}{d \rho^{2}}+\frac{l(l+1)}{\rho^{2}}\right) U_{l}(\rho)=U_{l}(\rho) \tag{14.16}
\end{equation*}
$$

We would not have done this before with the hydrogen atom since we had only finite $E=\hbar^{2} k^{2} / 2 \mu$. Now this can be anything.

Making one final substitution, $U_{l}(\rho)=\rho f_{l}(\rho)$ we can rewrite eq. (14.16) as

$$
\begin{equation*}
\left(\rho^{2} \frac{d^{2}}{d \rho^{2}}+2 \rho \frac{d}{d \rho}+\left(\rho^{2}-l(l+1)\right)\right) f_{l}=0 \tag{14.17}
\end{equation*}
$$

This is the spherical Bessel equation of order $l$ and has solutions called the Bessel and Neumann functions of order $l$, which are

$$
\begin{align*}
& j_{l}(\rho)=(-\rho)^{l}\left(\frac{1}{\rho} \frac{d}{d \rho}\right)^{l}\left(\frac{\sin \rho}{\rho}\right)  \tag{14.18a}\\
& n_{l}(\rho)=(-\rho)^{l}\left(\frac{1}{\rho} \frac{d}{d \rho}\right)^{l}\left(-\frac{\cos \rho}{\rho}\right) . \tag{14.18b}
\end{align*}
$$

We can easily calculate

$$
\begin{align*}
& U_{0}(\rho)=\rho j_{0}(\rho)=\sin \rho \\
& U_{1}(\rho)=\rho j_{1}(\rho)=-\cos \rho+\frac{\sin \rho}{\rho} \tag{14.19}
\end{align*}
$$

and can plug these into eq. (14.16) to verify that they are a solution. A more general proof looks a bit trickier. Observe that the Neumann functions are less well behaved at the origin. To calculate the first few Bessel and Neumann functions we first compute

$$
\begin{align*}
& \frac{1}{\rho} \frac{d}{d \rho} \frac{\sin \rho}{\rho}=\frac{1}{\rho}\left(\frac{\cos \rho}{\rho}-\frac{\sin \rho}{\rho^{2}}\right)  \tag{14.20}\\
&=\frac{\cos \rho}{\rho^{2}}-\frac{\sin \rho}{\rho^{3}} \\
&\left(\begin{array}{l}
\left.\frac{1}{\rho} \frac{d}{d \rho}\right)^{2} \frac{\sin \rho}{\rho}
\end{array}=\frac{1}{\rho}\left(-\frac{\sin \rho}{\rho^{2}}-2 \frac{\cos \rho}{\rho^{3}}-\frac{\cos \rho}{\rho^{3}}+3 \frac{\sin \rho}{\rho^{4}}\right)\right.  \tag{14.21}\\
&=\sin \rho\left(-\frac{1}{\rho^{3}}+\frac{3}{\rho^{5}}\right)-3 \frac{\cos \rho}{\rho^{4}},
\end{align*}
$$

and

$$
\begin{align*}
\frac{1}{\rho} \frac{d}{d \rho}-\frac{\cos \rho}{\rho} & =\frac{1}{\rho}\left(\frac{\sin \rho}{\rho}+\frac{\cos \rho}{\rho^{2}}\right)  \tag{14.22}\\
& =\frac{\sin \rho}{\rho^{2}}+\frac{\cos \rho}{\rho^{3}} \\
\left(\frac{1}{\rho} \frac{d}{d \rho}\right)^{2}-\frac{\cos \rho}{\rho} & =\frac{1}{\rho}\left(\frac{\cos \rho}{\rho^{2}}-2 \frac{\sin \rho}{\rho^{3}}-\frac{\sin \rho}{\rho^{3}}-3 \frac{\cos \rho}{\rho^{4}}\right)  \tag{14.23}\\
& =\cos \rho\left(\frac{1}{\rho^{3}}-\frac{3}{\rho^{5}}\right)-3 \frac{\sin \rho}{\rho^{4}},
\end{align*}
$$

so we find

$$
\begin{align*}
& j_{0}(\rho)=\frac{\sin \rho}{\rho} \\
& j_{1}(\rho)=\frac{\sin \rho}{\rho^{2}}-\frac{\cos \rho}{\rho}  \tag{14.24}\\
& j_{2}(\rho)=\sin \rho\left(-\frac{1}{\rho}+\frac{3}{\rho^{3}}\right)+\cos \rho\left(-\frac{3}{\rho^{2}}\right),
\end{align*}
$$

and

$$
\begin{align*}
& n_{0}(\rho)=-\frac{\cos \rho}{\rho} \\
& n_{1}(\rho)=-\frac{\cos \rho}{\rho^{2}}-\frac{\sin \rho}{\rho}  \tag{14.25}\\
& n_{2}(\rho)=\cos \rho\left(\frac{1}{\rho}-\frac{3}{\rho^{3}}\right)+\sin \rho\left(-\frac{3}{\rho^{2}}\right)
\end{align*}
$$

Observe that our radial functions $R(r)$ are proportional to these Bessel and Neumann functions

$$
\begin{align*}
R(r) & =\frac{u(r)}{r} \\
& =\frac{U(k r)}{r} \\
& =\left\{\begin{array}{l}
\frac{j_{l}(\rho) \rho}{r} \\
\frac{n_{l}(\rho) \rho}{r}
\end{array}\right.  \tag{14.26}\\
& =\left\{\begin{array}{l}
\frac{j_{l}(\rho) k l^{\prime}}{r} \\
\frac{n_{l}(\rho) k \gamma^{\prime}}{r}
\end{array},\right.
\end{align*}
$$

or

$$
\begin{equation*}
R(r) \sim j_{l}(\rho), n_{l}(\rho) \tag{14.27}
\end{equation*}
$$

### 14.4 LIMITS OF SPHERICAL BESSEL AND NEUMANN FUNCTIONS.

With $n$ ! ! denoting the double factorial, like factorial but skipping every other term

$$
\begin{equation*}
n!!=n(n-2)(n-4) \cdots, \tag{14.28}
\end{equation*}
$$

we can show that in the limit as $\rho \rightarrow 0$ we have

$$
\begin{equation*}
j_{l}(\rho) \rightarrow \frac{\rho^{l}}{(2 l+1)!!} \tag{14.29a}
\end{equation*}
$$

$$
\begin{equation*}
n_{l}(\rho) \rightarrow-\frac{(2 l-1)!!}{\rho^{(l+1)}} \tag{14.29b}
\end{equation*}
$$

(for the $l=0$ case, note that $(-1)!!=1$ by definition). Comparing this to our explicit expansion for $j_{1}(\rho)$ in eq. (14.24) where we appear to have a
$1 / \rho$ dependence for small $\rho$ it is not obvious that this would be the case. To compute this we need to start with a power series expansion for $\sin \rho / \rho$, which is well behaved at $\rho=0$ and then the result follows (done later). It is apparently also possible to show that as $\rho \rightarrow \infty$ we have

$$
\begin{align*}
& j_{l}(\rho) \rightarrow \frac{1}{\rho} \sin \left(\rho-\frac{l \pi}{2}\right)  \tag{14.30a}\\
& n_{l}(\rho) \rightarrow-\frac{1}{\rho} \cos \left(\rho-\frac{l \pi}{2}\right) . \tag{14.30b}
\end{align*}
$$

## 14.5 back to our problem.

For $r>r_{0}$ we can construct (for fixed $k$ ) a superposition of the spherical functions

$$
\begin{equation*}
\sum_{l} \sum_{m}\left(A_{l} j_{l}(k r)+B_{l} n_{l}(k r)\right) Y_{l}^{m}(\theta, \phi), \tag{14.31}
\end{equation*}
$$

we want outgoing waves, and as $r \rightarrow \infty$, we have

$$
\begin{align*}
& j_{l}(k r) \rightarrow \frac{\sin \left(k r-\frac{l \pi}{2}\right)}{k r}  \tag{14.32a}\\
& n_{l}(k r) \rightarrow-\frac{\cos \left(k r-\frac{l \pi}{2}\right)}{k r} . \tag{14.32b}
\end{align*}
$$

Put $A_{l} / B_{l}=-i$ for a given $l$ we have

$$
\begin{equation*}
\frac{1}{k r}\left(-i \frac{\sin \left(k r-\frac{l \pi}{2}\right)}{k r}-\frac{\cos \left(k r-\frac{l \pi}{2}\right)}{k r}\right) \sim \frac{1}{k r} e^{i(k r-\pi l / 2)} \tag{14.33}
\end{equation*}
$$

for

$$
\begin{equation*}
\sum_{l} \sum_{m} B_{l} \frac{1}{k r} e^{i(k r-\pi l / 2)} Y_{l}^{m}(\theta, \phi) . \tag{14.34}
\end{equation*}
$$

Making this choice to achieve outgoing waves (and factoring a $(-i)^{l}$ out of $B_{l}$ for some reason, we have another wave function that satisfies our Hamiltonian equation

$$
\begin{equation*}
\frac{e^{i k r}}{k r} \sum_{l} \sum_{m}(-1)^{l} B_{l} Y_{l}^{m}(\theta, \phi) \tag{14.35}
\end{equation*}
$$

The $B_{l}$ coefficients will depend on $V(r)$ for the incident wave $e^{i \mathbf{k} \cdot \mathbf{r}}$. Suppose we encapsulate that dependence in a helper function $f_{\mathbf{k}}(\theta, \phi)$ and write

$$
\begin{equation*}
\frac{e^{i k r}}{r} f_{\mathbf{k}}(\theta, \phi) \tag{14.36}
\end{equation*}
$$

We seek a solution $\psi_{\mathbf{k}}(\mathbf{r})$

$$
\begin{equation*}
\left(-\frac{\hbar^{2}}{2 \mu} \nabla^{2}+V(\mathbf{r})\right) \psi_{\mathbf{k}}(\mathbf{r})=\frac{\hbar^{2} \mathbf{k}^{2}}{2 \mu} \psi_{\mathbf{k}}(\mathbf{r}) \tag{14.37}
\end{equation*}
$$

where as $r \rightarrow \infty$

$$
\begin{equation*}
\psi_{\mathbf{k}}(\mathbf{r}) \rightarrow e^{i \mathbf{k} \cdot \mathbf{r}}+\frac{e^{i k r}}{r} f_{\mathbf{k}}(\theta, \phi) \tag{14.38}
\end{equation*}
$$

Note that for $r<r_{0}$ in general for finite $r, \psi_{k}(\mathbf{r})$, is much more complicated. This is the analogue of the plane wave result

$$
\begin{equation*}
\psi(x)=e^{i k x}+\beta_{k} e^{-i k x} . \tag{14.39}
\end{equation*}
$$

### 14.6 SCattering geometry and nomenclature.

We can think classically first, and imagine a scattering of a stream of particles barraging a target as in fig. 14.3. Here we assume that $d \Omega$ is far


Figure 14.3: Scattering cross section.
enough away that it includes no non-scattering particles. Write $P$ for the number density

$$
\begin{equation*}
P=\frac{\text { number of particles }}{\text { unit volume }} \tag{14.40}
\end{equation*}
$$

and

$$
\begin{equation*}
J=P v_{0}=\frac{\text { Number of particles flowing through }}{\text { a unit area in unit time }} . \tag{14.41}
\end{equation*}
$$

We want to count the rate of particles per unit time $d N$ through this solid angle $d \Omega$ and write

$$
\begin{equation*}
d N=J\left(\frac{d \sigma(\Omega)}{d \Omega}\right) d \Omega \tag{14.42}
\end{equation*}
$$

The factor

$$
\begin{equation*}
\frac{d \sigma(\Omega)}{d \Omega} \tag{14.43}
\end{equation*}
$$

is called the differential cross section, and has "units" of

$$
\begin{equation*}
\frac{\text { area }}{\text { steradians }} \tag{14.44}
\end{equation*}
$$

(recalling that steradians are radian like measures of solid angle [21]). The total number of particles through the volume per unit time is then

$$
\begin{equation*}
\int J \frac{d \sigma(\Omega)}{d \Omega} d \Omega=J \int \frac{d \sigma(\Omega)}{d \Omega} d \Omega=J \sigma \tag{14.45}
\end{equation*}
$$

where $\sigma$ is the total cross section and has units of area. The cross section $\sigma$ his the effective size of the area required to collect all particles, and characterizes the scattering, but is not necessarily entirely geometrical. For example, in photon scattering we may have frequency matching with atomic resonance, finding $\sigma \sim \lambda^{2}$, something that can be much bigger than the actual total area involved.

### 14.7 APPENDIX.

Q: Are Bessel and Neumann functions orthogonal?

Answer: There is an orthogonality relation, but it is not one of plain old multiplication. Curious about this, I find an orthogonality condition in [19]

$$
\begin{equation*}
\int_{0}^{\infty} J_{\alpha}(z) J_{\beta}(z) \frac{d z}{z}=\frac{2}{\pi} \frac{\sin \left(\frac{\pi}{2}(\alpha-\beta)\right)}{\alpha^{2}-\beta^{2}} \tag{14.46}
\end{equation*}
$$

from which we find for the spherical Bessel functions

$$
\begin{equation*}
\int_{0}^{\infty} j_{l}(\rho) j_{m}(\rho) d \rho=\frac{\sin \left(\frac{\pi}{2}(l-m)\right)}{(l+1 / 2)^{2}-(m+1 / 2)^{2}} \tag{14.47}
\end{equation*}
$$

Is this a satisfactory orthogonality integral? At a glance it does not appear to be well behaved for $l=m$, but perhaps the limit can be taken?

Deriving the large limit Bessel and Neumann function approximations. For eq. (14.30) we are referred to any "good book on electromagnetism" for details. I thought that perhaps the weighty [10] would be to be such a book, but it also leaves out the details. In $\S 16.1$ the spherical Bessel and Neumann functions are related to the plain old Bessel functions with

$$
\begin{align*}
& j_{l}(x)=\sqrt{\frac{\pi}{2 x}} J_{l+1 / 2}(x)  \tag{14.48a}\\
& n_{l}(x)=\sqrt{\frac{\pi}{2 x}} N_{l+1 / 2}(x) . \tag{14.48b}
\end{align*}
$$

Referring back to $\S 3.7$ of that text where the limiting forms of the Bessel functions are given

$$
\begin{align*}
& J_{v}(x) \rightarrow \sqrt{\frac{2}{\pi x}} \cos \left(x-\frac{v \pi}{2}-\frac{\pi}{4}\right)  \tag{14.49a}\\
& N_{v}(x) \rightarrow \sqrt{\frac{2}{\pi x}} \sin \left(x-\frac{v \pi}{2}-\frac{\pi}{4}\right) . \tag{14.49b}
\end{align*}
$$

This does give us our desired identities, but there is no hint in the text how one would derive eq. (14.49) from the power series that was computed by solving the Bessel equation.

Deriving the small limit Bessel and Neumann function approximations. Writing the sinc function in series form

$$
\begin{equation*}
\frac{\sin x}{x}=\sum_{k=0}^{\infty}(-1)^{k} \frac{x^{2 k}}{(2 k+1)!} \tag{14.50}
\end{equation*}
$$

we can differentiate easily

$$
\begin{align*}
\frac{1}{x} \frac{d}{d x} \frac{\sin x}{x} & =\sum_{k=1}^{\infty}(-1)^{k}(2 k) \frac{x^{2 k-2}}{(2 k+1)!} \\
& =(-1) \sum_{k=0}^{\infty}(-1)^{k}(2 k+2) \frac{x^{2 k}}{(2 k+3)!}  \tag{14.51}\\
& =(-1) \sum_{k=0}^{\infty}(-1)^{k} \frac{1}{2 k+3} \frac{x^{2 k}}{(2 k+1)!} .
\end{align*}
$$

Performing the derivative operation a second time we find

$$
\begin{align*}
\left(\frac{1}{x} \frac{d}{d x}\right)^{2} \frac{\sin x}{x} & =(-1) \sum_{k=1}^{\infty}(-1)^{k} \frac{1}{2 k+3}(2 k) \frac{x^{2 k-2}}{(2 k+1)!}  \tag{14.52}\\
& =\sum_{k=0}^{\infty}(-1)^{k} \frac{1}{2 k+5} \frac{1}{2 k+3} \frac{x^{2 k}}{(2 k+1)!} .
\end{align*}
$$

It appears reasonable to form the inductive hypotheses

$$
\begin{equation*}
\left(\frac{1}{x} \frac{d}{d x}\right)^{l} \frac{\sin x}{x}=(-1)^{l} \sum_{k=0}^{\infty}(-1)^{k} \frac{(2 k+1)!!}{(2(k+l)+1)!!} \frac{x^{2 k}}{(2 k+1)!}, \tag{14.53}
\end{equation*}
$$

and this proves to be correct. We find then that the spherical Bessel function has the power series expansion of

$$
\begin{equation*}
j_{l}(x)=\sum_{k=0}^{\infty}(-1)^{k} \frac{(2 k+1)!!}{(2(k+l)+1)!!} \frac{x^{2 k+l}}{(2 k+1)!}, \tag{14.54}
\end{equation*}
$$

and from this the Bessel function limit of eq. (14.29a) follows immediately. Finding the matching induction series for the Neumann functions is a bit harder. It is not really any more difficult to write it, but it is harder to put it in a tidy form that is. We find

$$
\begin{align*}
-\frac{\cos x}{x} & =-\sum_{k=0}^{\infty}(-1)^{k} \frac{x^{2 k-1}}{(2 k)!} \\
\frac{1}{x} \frac{d}{d x}-\frac{\cos x}{x} & =-\sum_{k=0}^{\infty}(-1)^{k} \frac{2 k-1}{2 k} \frac{x^{2 k-3}}{(2 k-2)!}  \tag{14.55}\\
\left(\frac{1}{x} \frac{d}{d x}\right)^{2}-\frac{\cos x}{x} & =-\sum_{k=0}^{\infty}(-1)^{k} \frac{(2 k-1)(2 k-3)}{2 k(2 k-2)} \frac{x^{2 k-3}}{(2 k-4)!} .
\end{align*}
$$

The general expression, after a bit of messing around (and I got it wrong the first time), can be found to be

$$
\begin{align*}
\left(\frac{1}{x} \frac{d}{d x}\right)^{l}-\frac{\cos x}{x}= & (-1)^{l+1} \sum_{k=0}^{l-1} \prod_{j=0}^{l-1}|2(k-j)-1| \frac{x^{2(k-l)-1}}{(2 k)!}  \tag{14.56}\\
& +(-1)^{l+1} \sum_{k=0}^{\infty}(-1)^{k} \frac{(2(k+l)-1)!!}{(2 k-1)!!} \frac{x^{2 k-1}}{(2(k+l)!}
\end{align*}
$$

We really only need the lowest order term (which dominates for small $x$ ) to confirm the small limit eq. (14.29b) of the Neumann function, and this follows immediately. For completeness, we note that the series expansion of the Neumann function is

$$
\begin{align*}
n_{l}(x)= & -\sum_{k=0}^{l-1} \prod_{j=0}^{l-1}|2(k-j)-1| \frac{x^{2 k-l-1}}{(2 k)!}  \tag{14.57}\\
& -\sum_{k=0}^{\infty}(-1)^{k} \frac{(2 k+3 l-1)!!}{(2 k-1)!!} \frac{x^{2 k-1}}{(2(k+l)!}
\end{align*}
$$

### 14.8 VERIFYING THE SOLUTION TO THE SPHERICAL BESSEL EQUATION.

One way to verify that eq. (14.18a) is a solution to the Bessel equation eq. (14.17) as claimed should be to substitute the series expression and verify that we get zero. Another way is to solve this equation directly. We have a regular singular point at the origin, so we look for solutions of the form

$$
\begin{equation*}
f=x^{r} \sum_{k=0}^{\infty} a_{k} x^{k} \tag{14.58}
\end{equation*}
$$

Writing our differential operator as

$$
\begin{equation*}
L=x^{2} \frac{d^{2}}{d x^{2}}+2 x \frac{d}{d x}+x^{2}-l(l+1) \tag{14.59}
\end{equation*}
$$

we get

$$
\begin{align*}
0 & =L f \\
& =\sum_{k=0}^{\infty} a_{k}((k+r)(k+r-1)+2(k+r)-l(l+1)) x^{k+r}+a_{k} x^{k+r+2} \\
& =a_{0}(r(r+1)-l(l+1)) x^{r}  \tag{14.60}\\
& +a_{1}((r+1)(r+2)-l(l+1)) x^{r+1} \\
& +\sum_{k=2}^{\infty} a_{k}\left((k+r)(k+r-1)+2(k+r)-l(l+1)+a_{k-2}\right) x^{k+r} .
\end{align*}
$$

Since we require this to be zero for all $x$ including non-zero values, we must have constraints on $r$. Assuming first that $a_{0}$ is non-zero we must then have

$$
\begin{equation*}
0=r(r+1)-l(l+1) \tag{14.61}
\end{equation*}
$$

One solution is obviously $r=l$. Assuming we have another solution $r=l+k$ for some integer $k$ we find that $r=-l-1$ is also a solution. Restricting attention first to $r=l$, we must have $a_{1}=0$ since for nonnegative $l$ we have $(l+1)(l+2)-l(l+1)=2(l+1) \neq 0$. Thus for non-zero $a_{0}$ we find that our function is of the form

$$
\begin{equation*}
f=\sum_{k} a_{2 k} x^{2 k+l} \tag{14.62}
\end{equation*}
$$

It does not matter that we started with $a_{0} \neq 0$. If we instead start with $a_{1} \neq 0$ we find that we must have $r=l-1,-l-2$, so end up with exactly the same functional form as eq. (14.62). It ends up slightly simpler if we start with eq. (14.62) instead, since we now know that we do not have any odd powered $a_{k}$ 's to deal with. Doing so we find

$$
\begin{align*}
0 & =L f \\
& =\sum_{k=0}^{\infty} a_{2 k}((2 k+l)(2 k+l-1)+2(2 k+l)-l(l+1)) x^{2 k+l}+a_{2 k} x^{2 k+l+2} \\
& =\sum_{k=1}^{\infty}\left(a_{2 k} 2 k(2(k+l)+1)+a_{2(k-1)}\right) x^{2 k+l} . \tag{14.63}
\end{align*}
$$

We find

$$
\begin{equation*}
\frac{a_{2 k}}{a_{2(k-1)}}=\frac{-1}{2 k(2(k+l)+1)} \tag{14.64}
\end{equation*}
$$

Proceeding recursively, we find

$$
\begin{equation*}
f=a_{0}(2 l+1)!!\sum_{k=0}^{\infty} \frac{(-1)^{k}}{(2 k)!!(2(k+l)+1)!!} x^{2 k+l} \tag{14.65}
\end{equation*}
$$

With $a_{0}=1 /(2 l+1)!$ ! and the observation that

$$
\begin{equation*}
\frac{1}{(2 k)!!}=\frac{(2 k+1)!!}{(2 k+1)!} \tag{14.66}
\end{equation*}
$$

we have $f=j_{l}(x)$ as given in eq. (14.54).
If we do the same for the $r=-l-1$ case, we find

$$
\begin{equation*}
\frac{a_{2 k}}{a_{2(k-1)}}=\frac{-1}{2 k(2(k-l)-1)}, \tag{14.67}
\end{equation*}
$$

and find

$$
\begin{equation*}
\frac{a_{2 k}}{a_{0}}=\frac{(-1)^{k}}{(2 k)!!(2(k-l)-1)(2(k-l)-3) \cdots(-2 l+1)} . \tag{14.68}
\end{equation*}
$$

Flipping signs around, we can rewrite this as

$$
\begin{equation*}
\frac{a_{2 k}}{a_{0}}=\frac{1}{(2 k)!!(2(l-k)+1)(2(l-k)+3) \cdots(2 l-1)} \tag{14.69}
\end{equation*}
$$

For those values of $l>k$ we can write this as

$$
\begin{equation*}
\frac{a_{2 k}}{a_{0}}=\frac{(2(l-k)-1)!!}{(2 k)!!(2 l-1)!!} \tag{14.70}
\end{equation*}
$$

Comparing to the small limit eq. (14.29b), the $k=0$ term, we find that we must have

$$
\begin{equation*}
\frac{a_{0}}{(2 l-1)!!}=-1 . \tag{14.71}
\end{equation*}
$$

After some play we find

$$
a_{2 k}=\left\{\begin{array}{ll}
-\frac{(2(l-k)-1)!!}{(2 k)!!} & \text { if } l \geq k  \tag{14.72}\\
\frac{(-1)^{k-l+1}}{(2 k)!(2(k-l)-1)!!} & \text { if } l \leq k
\end{array} .\right.
$$

Putting this all together we have

$$
\begin{equation*}
n_{l}(x)=-\sum_{0 \leq k \leq l}(2(l-k)-1)!!\frac{x^{2 k-l-1}}{(2 k)!!}-\sum_{l<k} \frac{(-1)^{k-l}}{(2(k-l)-1)!!} \frac{x^{2 k-l-1}}{(2 k)!!} \tag{14.73}
\end{equation*}
$$

### 14.9 SCATTERING CROSS SECTIONS.

## READING: §20 [4]

Recall that we are studying the case of a potential that is zero outside of a fixed bound, $V(\mathbf{r})=0$ for $r>r_{0}$, as in fig. 14.4, and were looking for


Figure 14.4: Bounded potential.
solutions to Schrödinger's equation

$$
\begin{equation*}
-\frac{\hbar^{2}}{2 \mu} \nabla^{2} \psi_{\mathbf{k}}(\mathbf{r})+V(\mathbf{r}) \psi_{\mathbf{k}}(\mathbf{r})=\frac{\hbar^{2} \mathbf{k}^{2}}{2 \mu} \psi_{\mathbf{k}}(\mathbf{r}) \tag{14.74}
\end{equation*}
$$

in regions of space, where $r>r_{0}$ is very large. We found

$$
\begin{equation*}
\psi_{\mathbf{k}}(\mathbf{r}) \sim e^{i \mathbf{k} \cdot \mathbf{r}}+\frac{e^{i k r}}{r} f_{\mathbf{k}}(\theta, \phi) \tag{14.75}
\end{equation*}
$$

For $r \leq r_{0}$ this will be something much more complicated. To study scattering we will use the concept of probability flux as in electromagnetism

$$
\begin{equation*}
\boldsymbol{\nabla} \cdot \mathbf{j}+\dot{\rho}=0 \tag{14.76}
\end{equation*}
$$

Using

$$
\begin{equation*}
\psi(\mathbf{r}, t)=\psi_{\mathbf{k}}(\mathbf{r})^{*} \psi_{\mathbf{k}}(\mathbf{r}) \tag{14.77}
\end{equation*}
$$

we find

$$
\begin{equation*}
\mathbf{j}(\mathbf{r}, t)=\frac{\hbar}{2 \mu i}\left(\psi_{\mathbf{k}}(\mathbf{r})^{*} \boldsymbol{\nabla} \psi_{\mathbf{k}}(\mathbf{r})-\left(\boldsymbol{\nabla} \psi_{\mathbf{k}}^{*}(\mathbf{r})\right) \psi_{\mathbf{k}}(\mathbf{r})\right) \tag{14.78}
\end{equation*}
$$

when

$$
\begin{equation*}
-\frac{\hbar^{2}}{2 \mu} \nabla^{2} \psi_{\mathbf{k}}(\mathbf{r})+V(\mathbf{r}) \psi_{\mathbf{k}}(\mathbf{r})=i \hbar \frac{\partial \psi_{\mathbf{k}}(\mathbf{r})}{\partial t} \tag{14.79}
\end{equation*}
$$

In a fashion similar to what we did in the 1D case, let us suppose that we can write our wave function

$$
\begin{equation*}
\psi\left(\mathbf{r}, t_{\text {initial }}\right)=\int d^{3} k \alpha\left(\mathbf{k}, t_{\text {initial }}\right) \psi_{\mathbf{k}}(\mathbf{r}) \tag{14.80}
\end{equation*}
$$

and treat the scattering as the scattering of a plane wave front (idealizing a set of wave packets) off of the object of interest as depicted in fig. 14.5. We assume that our incoming particles are sufficiently localized in $k$ space


Figure 14.5: plane wave front incident on particle.
as depicted in the idealized representation of fig. 14.6. We assume that


Figure 14.6: k space localized wave packet.
$\alpha\left(\mathbf{k}, t_{\text {initial }}\right)$ is localized.

$$
\begin{equation*}
\psi\left(\mathbf{r}, t_{\text {initial }}\right)=\int d^{3} k\left(\alpha\left(\mathbf{k}, t_{\text {initial }}\right) e^{i k_{z} z}+\alpha\left(\mathbf{k}, t_{\text {initial }}\right) \frac{e^{i k r}}{r} f_{\mathbf{k}}(\theta, \phi)\right) \tag{14.81}
\end{equation*}
$$

We suppose that

$$
\begin{equation*}
\alpha\left(\mathbf{k}, t_{\text {initial }}\right)=\alpha(\mathbf{k}) e^{-i \hbar k^{2} t_{\text {initial }} / 2 \mu} \tag{14.82}
\end{equation*}
$$

where this is chosen $\left(\alpha\left(\mathbf{k}, t_{\text {initial }}\right)\right.$ is built in this fashion) so that this is non-zero for $z$ large in magnitude and negative. This last integral can be approximated

$$
\begin{align*}
\int d^{3} k \alpha\left(\mathbf{k}, t_{\text {initial }}\right) \frac{e^{i k r}}{r} f_{\mathbf{k}}(\theta, \phi) & \approx \frac{f_{\mathbf{k}_{0}}(\theta, \phi)}{r} \int d^{3} k \alpha\left(\mathbf{k}, t_{\text {initial }}\right) e^{i k r}  \tag{14.83}\\
& \rightarrow 0
\end{align*}
$$

This is very much like the 1D case where we found no reflected component for our initial time. We will normally look in a locality well away from the wave front as indicted in fig. 14.7. There are situations where we


Figure 14.7: point of measurement of scattering cross section.
do look in the locality of the wave front that has been scattered.
Incoming wave. Our income wave is of the form

$$
\begin{equation*}
\psi_{i}=A e^{i k z} e^{-i \hbar k^{2} t / 2 \mu} \tag{14.84}
\end{equation*}
$$

Here we have made the approximation that $k=|\mathbf{k}| \sim k_{z}$. We can calculate the probability current

$$
\begin{equation*}
\mathbf{j}=\hat{\mathbf{z}} \frac{\hbar k}{\mu} A, \tag{14.85}
\end{equation*}
$$

(notice the $v=p / m$ like term above, with $p=\hbar k$ ). For the scattered wave (dropping $A$ factor)

$$
\begin{aligned}
\mathbf{j} & =\frac{\hbar}{2 \mu i}\left(f_{\mathbf{k}}^{*}(\theta, \phi) \frac{e^{-i k r}}{r} \nabla\left(f_{\mathbf{k}}(\theta, \phi) \frac{e^{i k r}}{r}\right)-\nabla\left(f_{\mathbf{k}}^{*}(\theta, \phi) \frac{e^{-i k r}}{r}\right) f_{\mathbf{k}}(\theta, \phi) \frac{e^{i k r}}{r}\right) \\
& \approx \frac{\hbar}{2 \mu i}\left(f_{\mathbf{k}}^{*}(\theta, \phi) \frac{e^{-i k r}}{r} i k \hat{\mathbf{r}} f_{\mathbf{k}}(\theta, \phi) \frac{e^{i k r}}{r}-f_{\mathbf{k}}^{*}(\theta, \phi) \frac{e^{-i k r}}{r}(-i k \hat{\mathbf{r}}) f_{\mathbf{k}}(\theta, \phi) \frac{e^{i k r}}{r}\right),
\end{aligned}
$$

We find that the radial portion of the current density is

$$
\begin{align*}
\hat{\mathbf{r}} \cdot \mathbf{j} & =\frac{\hbar}{2 \mu i}|f|^{2} \frac{2 i k}{r^{2}} \\
& =\frac{\hbar k}{\mu} \frac{1}{r^{2}}|f|^{2} \tag{14.87}
\end{align*}
$$

and the flux through our element of solid angle is

$$
\begin{align*}
\hat{\mathbf{r}} d A \cdot \mathbf{j} & =\frac{\text { probability }}{\text { unit area per time }} \times \text { area } \\
& =\frac{\text { probability }}{\text { unit time }} \\
& =\frac{\hbar k}{\mu} \frac{\left|f_{\mathbf{k}}(\theta, \phi)\right|^{2}}{r^{2}} r^{2} d \Omega  \tag{14.88}\\
& =\frac{\hbar k}{\mu}\left|f_{\mathbf{k}}(\theta, \phi)\right|^{2} d \Omega \\
& d \sigma / d \Omega \\
& =j_{\text {incoming }}\left|f_{\mathbf{k}}(\theta, \phi)\right|^{2} d \Omega .
\end{align*}
$$

We identify the scattering cross section above

$$
\begin{align*}
& \frac{d \sigma}{d \Omega}=\left|f_{\mathbf{k}}(\theta, \phi)\right|^{2}  \tag{14.89}\\
& \sigma=\int\left|f_{\mathbf{k}}(\theta, \phi)\right|^{2} d \Omega \tag{14.90}
\end{align*}
$$

We have been somewhat unrealistic here since we have used a plane wave approximation. Should that plane wave be replaced by a more general wave packet (as crudely sketched in fig. 14.8), we will actually obtain the same answer. For details we are referred to [12] and [15].

Working towards a solution. We have done a bunch of stuff here but are not much closer to a real solution because we do not actually know what $f_{\mathbf{k}}$ is. Let us write Schrödinger

$$
\begin{equation*}
-\frac{\hbar^{2}}{2 \mu} \nabla^{2} \psi_{\mathbf{k}}(\mathbf{r})+V(\mathbf{r}) \psi_{\mathbf{k}}(\mathbf{r})=\frac{\hbar^{2} \mathbf{k}^{2}}{2 \mu} \psi_{\mathbf{k}}(\mathbf{r}) \tag{14.91}
\end{equation*}
$$



Figure 14.8: Plane wave vs packet wave front.
instead as

$$
\begin{equation*}
\left(\boldsymbol{\nabla}^{2}+\mathbf{k}^{2}\right) \psi_{\mathbf{k}}(\mathbf{r})=s(\mathbf{r}), \tag{14.92}
\end{equation*}
$$

where

$$
\begin{equation*}
s(\mathbf{r})=\frac{2 \mu}{\hbar^{2}} V(\mathbf{r}) \psi_{\mathbf{k}}(\mathbf{r}) \tag{14.93}
\end{equation*}
$$

where $s(\mathbf{r})$ is really the particular solution to this differential problem. We want

$$
\begin{equation*}
\psi_{\mathbf{k}}(\mathbf{r})=\psi_{\mathbf{k}}^{\text {homogeneous }}(\mathbf{r})+\psi_{\mathbf{k}}^{\text {particular }}(\mathbf{r}) \tag{14.94}
\end{equation*}
$$

and

$$
\begin{equation*}
\psi_{\mathbf{k}}^{\text {homogeneous }}(\mathbf{r})=e^{i \mathbf{k} \cdot \mathbf{r}} \tag{14.95}
\end{equation*}
$$

## READING: §20 [4]

We have been arguing that we can write the stationary equation

$$
\begin{equation*}
\left(\boldsymbol{\nabla}^{2}+\mathbf{k}^{2}\right) \psi_{\mathbf{k}}(\mathbf{r})=s(\mathbf{r}) \tag{15.1}
\end{equation*}
$$

with

$$
\begin{equation*}
s(\mathbf{r})=\frac{2 \mu}{\hbar^{2}} V(\mathbf{r}) \psi_{\mathbf{k}}(\mathbf{r}) \tag{15.2}
\end{equation*}
$$

$$
\begin{equation*}
\psi_{\mathbf{k}}(\mathbf{r})=\psi_{\mathbf{k}}^{\text {homogeneous }}(\mathbf{r})+\psi_{\mathbf{k}}^{\text {particular }}(\mathbf{r}) \tag{15.3}
\end{equation*}
$$

Introduce Green function

$$
\begin{equation*}
\left(\boldsymbol{\nabla}^{2}+\mathbf{k}^{2}\right) G^{0}\left(\mathbf{r}, \mathbf{r}^{\prime}\right)=\delta\left(\mathbf{r}-\mathbf{r}^{\prime}\right) \tag{15.4}
\end{equation*}
$$

Suppose that I can find $G^{0}\left(\mathbf{r}, \mathbf{r}^{\prime}\right)$, then

$$
\begin{equation*}
\psi_{\mathbf{k}}^{\text {particular }}(\mathbf{r})=\int G^{0}\left(\mathbf{r}, \mathbf{r}^{\prime}\right) s\left(\mathbf{r}^{\prime}\right) d^{3} \mathbf{r}^{\prime} \tag{15.5}
\end{equation*}
$$

It turns out that finding the Green's function $G^{0}\left(\mathbf{r}, \mathbf{r}^{\prime}\right)$ is not so hard. Note the following, for $k=0$, we have

$$
\begin{equation*}
\nabla^{2} G_{0}^{0}\left(\mathbf{r}, \mathbf{r}^{\prime}\right)=\delta\left(\mathbf{r}-\mathbf{r}^{\prime}\right) \tag{15.6}
\end{equation*}
$$

(where a zero subscript is used to mark the $k=0$ case). We know this Green's function from electrostatics, and conclude that

$$
\begin{equation*}
G_{0}^{0}\left(\mathbf{r}, \mathbf{r}^{\prime}\right)=-\frac{1}{4 \pi} \frac{1}{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|} \tag{15.7}
\end{equation*}
$$

For $\mathbf{r} \neq \mathbf{r}^{\prime}$ we can easily show that

$$
\begin{equation*}
G^{0}\left(\mathbf{r}, \mathbf{r}^{\prime}\right)=-\frac{1}{4 \pi} \frac{e^{i k\left|\mathbf{r}-\mathbf{r}^{\prime}\right|}}{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|} \tag{15.8}
\end{equation*}
$$

This is correct for all $\mathbf{r}$ because it also gives the right limit as $\mathbf{r} \rightarrow \mathbf{r}^{\prime}$. This argument was first given by Lorentz. An outline for a derivation, utilizing the usual Fourier transform and contour integration arguments for these Green's derivations, can be found in $\S 7.4$ of [3]. A direct verification, not quite as easy as claimed can be found in L. We can now write our particular solution

$$
\begin{equation*}
\psi_{\mathbf{k}}(\mathbf{r})=e^{i \mathbf{k} \cdot \mathbf{r}}-\frac{1}{4 \pi} \int \frac{e^{i k\left|\mathbf{r}-\mathbf{r}^{\prime}\right|}}{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|} s\left(\mathbf{r}^{\prime}\right) d^{3} \mathbf{r}^{\prime} \tag{15.9}
\end{equation*}
$$

This is of no immediate help since we do not know $\psi_{\mathbf{k}}(\mathbf{r})$ and that is embedded in $s(\mathbf{r})$.

$$
\begin{equation*}
\psi_{\mathbf{k}}(\mathbf{r})=e^{i \mathbf{k} \cdot \mathbf{r}}-\frac{2 \mu}{4 \pi \hbar^{2}} \int \frac{e^{i k\left|\mathbf{r}-\mathbf{r}^{\prime}\right|}}{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|} V\left(\mathbf{r}^{\prime}\right) \psi_{\mathbf{k}}\left(\mathbf{r}^{\prime}\right) d^{3} \mathbf{r}^{\prime} \tag{15.10}
\end{equation*}
$$

Now look at this for $\mathbf{r} \gg \mathbf{r}^{\prime}$

$$
\begin{align*}
\left|\mathbf{r}-\mathbf{r}^{\prime}\right| & =\left(\mathbf{r}^{2}+\left(\mathbf{r}^{\prime}\right)^{2}-2 \mathbf{r} \cdot \mathbf{r}^{\prime}\right)^{1 / 2} \\
& =r\left(1+\frac{\left(\mathbf{r}^{\prime}\right)^{2}}{\mathbf{r}^{2}}-2 \frac{1}{\mathbf{r}^{2}} \mathbf{r} \cdot \mathbf{r}^{\prime}\right)^{1 / 2} \\
& =r\left(1-\frac{1}{2} \frac{2}{\mathbf{r}^{2}} \mathbf{r} \cdot \mathbf{r}^{\prime}+O\left(\frac{r^{\prime}}{r}\right)^{2}\right)^{1 / 2}  \tag{15.11}\\
& =r-\hat{\mathbf{r}} \cdot \mathbf{r}^{\prime}+O\left(\frac{r^{\prime 2}}{r}\right)
\end{align*}
$$

We get

$$
\begin{align*}
\psi_{\mathbf{k}}(\mathbf{r}) & \rightarrow e^{i \mathbf{k} \cdot \mathbf{r}}-\frac{2 \mu}{4 \pi \hbar^{2}} \frac{e^{i k r}}{r} \int e^{-i k \hat{\mathbf{r}} \cdot \mathbf{r}^{\prime}} V\left(\mathbf{r}^{\prime}\right) \psi_{\mathbf{k}}\left(\mathbf{r}^{\prime}\right) d^{3} \mathbf{r}^{\prime}  \tag{15.12}\\
& =e^{i \mathbf{k} \cdot \mathbf{r}}+f_{\mathbf{k}}(\theta, \phi) \frac{e^{i k r}}{r}
\end{align*}
$$

where

$$
\begin{equation*}
f_{\mathbf{k}}(\theta, \phi)=-\frac{\mu}{2 \pi \hbar^{2}} \int e^{-i k \hat{\mathbf{r}} \cdot \mathbf{r}^{\prime}} V\left(\mathbf{r}^{\prime}\right) \psi_{\mathbf{k}}\left(\mathbf{r}^{\prime}\right) d^{3} \mathbf{r}^{\prime} \tag{15.13}
\end{equation*}
$$

If the scattering is weak we have the Born approximation

$$
\begin{equation*}
f_{\mathbf{k}}(\theta, \phi)=-\frac{\mu}{2 \pi \hbar^{2}} \int e^{-i k \hat{\mathbf{r}} \cdot \mathbf{r}^{\prime}} V\left(\mathbf{r}^{\prime}\right) e^{i \mathbf{k} \cdot \mathbf{r}^{\prime}} d^{3} \mathbf{r}^{\prime} \tag{15.14}
\end{equation*}
$$

or

$$
\begin{equation*}
\psi_{\mathbf{k}}(\mathbf{r})=e^{i \mathbf{k} \cdot \mathbf{r}}-\frac{\mu}{2 \pi \hbar^{2}} \frac{e^{i k r}}{r} \int e^{-i k \hat{\mathbf{r}} \cdot \mathbf{r}^{\prime}} V\left(\mathbf{r}^{\prime}\right) e^{i \mathbf{k} \cdot \mathbf{r}^{\prime}} d^{3} \mathbf{r}^{\prime} \tag{15.15}
\end{equation*}
$$

Should we wish to make a further approximation, we can take the wave function resulting from application of the Born approximation, and use that a second time. This gives us the "Born again" approximation of

$$
\begin{align*}
\psi_{\mathbf{k}}(\mathbf{r})= & e^{i \mathbf{k} \cdot \mathbf{r}}-\frac{\mu}{2 \pi \hbar^{2}} \frac{e^{i k r}}{r} \int e^{-i k \hat{\mathbf{r}} \cdot \mathbf{r}^{\prime}} V\left(\mathbf{r}^{\prime}\right) \times \\
& \left(e^{i \mathbf{k} \cdot \mathbf{r}^{\prime}}-\frac{\mu}{2 \pi \hbar^{2}} \frac{e^{i k r^{\prime}}}{r^{\prime}} \int e^{-i k \hat{\mathbf{r}}^{\prime} \cdot \mathbf{r}^{\prime \prime}} V\left(\mathbf{r}^{\prime \prime}\right) e^{i \mathbf{k} \cdot \mathbf{r}^{\prime \prime}} d^{3} \mathbf{r}^{\prime \prime}\right) d^{3} \mathbf{r}^{\prime} \\
= & e^{i \mathbf{k} \cdot \mathbf{r}}-\frac{\mu}{2 \pi \hbar^{2}} \frac{e^{i k r}}{r} \int e^{-i k \hat{\mathbf{r}} \cdot \mathbf{r}^{\prime}} V\left(\mathbf{r}^{\prime}\right) e^{i \mathbf{k} \cdot \mathbf{r}^{\prime}} d^{3} \mathbf{r}^{\prime} \\
& +\frac{\mu^{2}}{(2 \pi)^{2} \hbar^{4}} \frac{e^{i k r}}{r} \int e^{-i k \hat{\mathbf{r}} \cdot \mathbf{r}^{\prime}} V\left(\mathbf{r}^{\prime}\right) \frac{e^{i k r^{\prime}}}{r^{\prime}} \int e^{-i k \mathbf{r}^{\prime} \cdot \mathbf{r}^{\prime \prime}} V\left(\mathbf{r}^{\prime \prime}\right) e^{i \mathbf{k} \cdot \mathbf{r}^{\prime \prime}} d^{3} \mathbf{r}^{\prime \prime} d^{3} \mathbf{r}^{\prime} \tag{15.16}
\end{align*}
$$

Consider

$$
\begin{equation*}
H_{0}=\frac{P^{2}}{2 m}+\frac{1}{2} m \omega^{2} X^{2} \tag{A.1}
\end{equation*}
$$

Since it has been a while let us compute the raising and lowering factorization that was used so extensively for this problem. It was of the form

$$
\begin{equation*}
H_{0}=(a X-i b P)(a X+i b P)+\cdots \tag{A.2}
\end{equation*}
$$

Why this factorization has an imaginary in it is a good question. It is not one that is given any sort of rationale in the text [4]. It is clear that we want $a=\sqrt{m / 2} \omega$ and $b=1 / \sqrt{2 m}$. The difference is then

$$
\begin{equation*}
H_{0}-(a X-i b P)(a X+i b P)=-i a b[X, P]=-i \frac{\omega}{2}[X, P] \tag{A.3}
\end{equation*}
$$

That commutator is an $i \hbar$ value, but what was the sign? Let us compute so we do not get it wrong

$$
\begin{align*}
{[x, p] \psi } & =-i \hbar\left[x, \partial_{x}\right] \psi \\
& =-i \hbar\left(x \partial_{x} \psi-\partial_{x}(x \psi)\right)  \tag{A.4}\\
& =-i \hbar(-\psi) \\
& =i \hbar \psi
\end{align*}
$$

so we have

$$
\begin{equation*}
H_{0}=\left(\omega \sqrt{\frac{m}{2}} X-i \sqrt{\frac{1}{2 m}} P\right)\left(\omega \sqrt{\frac{m}{2}} X+i \sqrt{\frac{1}{2 m}} P\right)+\frac{\hbar \omega}{2} . \tag{A.5}
\end{equation*}
$$

Factoring out an $\hbar \omega$ produces the form of the Hamiltonian that we used before

$$
\begin{equation*}
H_{0}=\hbar \omega\left(\left(\sqrt{\frac{m \omega}{2 \hbar}} X-i \sqrt{\frac{1}{2 m \hbar \omega}} P\right)\left(\sqrt{\frac{m \omega}{2 \hbar}} X+i \sqrt{\frac{1}{2 m \hbar \omega}} P\right)+\frac{1}{2}\right) \tag{A.6}
\end{equation*}
$$

The factors were labeled the uppering $\left(a^{\dagger}\right)$ and lowering (a) operators respectively, and written

$$
\begin{align*}
H_{0} & =\hbar \omega\left(a^{\dagger} a+\frac{1}{2}\right) \\
a & =\sqrt{\frac{m \omega}{2 \hbar}} X+i \sqrt{\frac{1}{2 m \hbar \omega}} P  \tag{A.7}\\
a^{\dagger} & =\sqrt{\frac{m \omega}{2 \hbar}} X-i \sqrt{\frac{1}{2 m \hbar \omega}} P .
\end{align*}
$$

Observe that we can find the inverse relations

$$
\begin{align*}
& X=\sqrt{\frac{\hbar}{2 m \omega}}\left(a+a^{\dagger}\right) \\
& P=i \sqrt{\frac{m \hbar \omega}{2}}\left(a^{\dagger}-a\right) . \tag{A.8}
\end{align*}
$$

Question: What is a good reason that we chose this particular factorization? For example, a quick computation shows that we could have also picked

$$
\begin{equation*}
H_{0}=\hbar \omega\left(a a^{\dagger}-\frac{1}{2}\right) \tag{A.9}
\end{equation*}
$$

I do not know that answer. That said, this second factorization is useful in that it provides the commutator relation between the raising and lowering operators, since subtracting eq. (A.9) and eq. (A.7) yields

$$
\begin{equation*}
\left[a, a^{\dagger}\right]=1 . \tag{A.10}
\end{equation*}
$$

If we suppose that we have eigenstates for the operator $a^{\dagger} a$ of the form

$$
\begin{equation*}
a^{\dagger} a|n\rangle=\lambda_{n}|n\rangle, \tag{A.11}
\end{equation*}
$$

then the problem of finding the eigensolution of $H_{0}$ reduces to solving this problem. Because $a^{\dagger} a$ commutes with $1 / 2$, an eigenstate of $a^{\dagger} a$ is also an eigenstate of $H_{0}$. Utilizing eq. (A.10) we then have

$$
\begin{align*}
a^{\dagger} a(a|n\rangle) & =\left(a a^{\dagger}-1\right) a|n\rangle \\
& =a\left(a^{\dagger} a-1\right)|n\rangle \\
& =a\left(\lambda_{n}-1\right)|n\rangle  \tag{A.12}\\
& =\left(\lambda_{n}-1\right) a|n\rangle
\end{align*}
$$

so we see that $a|n\rangle$ is an eigenstate of $a^{\dagger} a$ with eigenvalue $\lambda_{n}-1$. Similarly for the raising operator

$$
\begin{align*}
a^{\dagger} a\left(a^{\dagger}|n\rangle\right) & \left.=a^{\dagger}\left(a a^{\dagger}\right)|n\rangle\right) \\
& \left.=a^{\dagger}\left(a^{\dagger} a+1\right)|n\rangle\right)  \tag{A.13}\\
& \left.=a^{\dagger}\left(\lambda_{n}+1\right)|n\rangle\right)
\end{align*}
$$

and find that $a^{\dagger}|n\rangle$ is also an eigenstate of $a^{\dagger} a$ with eigenvalue $\lambda_{n}+1$.
Supposing that there is a lowest energy level (because the potential $V(x)=m \omega x^{2} / 2$ has a lower bound of zero) then the state $|0\rangle$ for which the energy is the lowest when operated on by $a$ we have

$$
\begin{equation*}
a|0\rangle=0 . \tag{A.14}
\end{equation*}
$$

Thus

$$
\begin{equation*}
a^{\dagger} a|0\rangle=0 \tag{A.15}
\end{equation*}
$$

and

$$
\begin{equation*}
\lambda_{0}=0 . \tag{A.16}
\end{equation*}
$$

This seems like a small bit of slight of hand, since it sneakily supplies an integer value to $\lambda_{0}$ where up to this point 0 was just a label. If the eigenvalue equation we are trying to solve for the Hamiltonian is

$$
\begin{equation*}
H_{0}|n\rangle=E_{n}|n\rangle, \tag{A.17}
\end{equation*}
$$

then we must have

$$
\begin{equation*}
E_{n}=\hbar \omega\left(\lambda_{n}+\frac{1}{2}\right)=\hbar \omega\left(n+\frac{1}{2}\right) . \tag{A.18}
\end{equation*}
$$

## a. 1 problems.

Exercise A. 1 Hermite recurrence relations.
Recall that Hermite polynomials are basis functions for the SHO wave functions ([4] eq. 9.58). Show that

$$
\frac{d}{d x} H_{n}(x)=2 n H_{n-1}(x) .
$$

## Answer for Exercise A. 1

We could just look up this formula, and if we did so can find it in various places, such as [1], or [6, Eq. 18.9.25].

To derive it ourself, we can start with the Rodgigues' formula definition of the Hermite polynomials

$$
\begin{equation*}
H_{n}=(-1)^{n} e^{x^{2}} \frac{d^{n}}{d x^{n}} e^{-x^{2}} \tag{A.19}
\end{equation*}
$$

Let us write $D=d / d x$, and take the derivative of $H_{n}$

$$
\begin{align*}
(-1)^{n} D H_{n} & =D\left(e^{x^{2}} D^{n} e^{-x^{2}}\right) \\
& =2 x e^{x^{2}} D^{n} e^{-x^{2}}+e^{x^{2}} D^{n}\left(-2 x e^{-x^{2}}\right) \\
& =2 x e^{x^{2}} D^{n} e^{-x^{2}}+e^{x^{2}} \sum_{k=0}^{n}\binom{n}{k} D^{k}(-2 x) D^{n-k} e^{-x^{2}} \\
& =2 x e^{x^{2}} D^{n} e^{-x^{2}}+e^{x^{2}} \sum_{k=0}^{1}\binom{n}{k} D^{k}(-2 x) D^{n-k} e^{-x^{2}}  \tag{A.20}\\
& =2 x e^{x^{2}} D^{n} e^{-x^{2}}+e^{x^{2}}\left(-2 x D^{n} e^{-x^{2}}-2 n D^{n-1} e^{-x^{2}}\right) \\
& =-2 n e^{x^{2}} D^{n-1} e^{-x^{2}} .
\end{align*}
$$

Rearranging proves the result.

On the quiz we were given a three state system $|1\rangle,|2\rangle$ and $|3\rangle$, and a two state system $|a\rangle,|b\rangle$, and were asked to show that the composite system can be entangled. I had trouble with this, having not seen any examples of this and subsequently filing away entanglement in the "abstract stuff that has no current known application" bit bucket, and then forgetting about it. Let us generate a concrete example of entanglement, and consider the very simplest direct product spaces. What is the simplest composite state that we can create? Suppose we have a pair of two state systems, say,

$$
\begin{align*}
& |1\rangle=\left[\begin{array}{l}
1 \\
0
\end{array}\right] \in H_{1}  \tag{B.1}\\
& |2\rangle=\left[\begin{array}{l}
0 \\
1
\end{array}\right] \in H_{1}
\end{align*}
$$

and

$$
\begin{align*}
& \langle x \mid+\rangle=\frac{e^{i k x}}{\sqrt{2 \pi}}, \text { where }|+\rangle \in H_{2} \\
& \langle x \mid+\rangle=\frac{e^{-i k x}}{\sqrt{2 \pi}} \text { where }|-\rangle \in H_{2} \tag{B.2}
\end{align*}
$$

We can now enumerate the space of possible operators

$$
\begin{align*}
A & \in a_{11++}|1\rangle\langle 1| \otimes|+\rangle\langle+|+a_{11+-}|1\rangle\langle 1| \otimes|+\rangle\langle-| \\
& +a_{11-+}|1\rangle\langle 1| \otimes|-\rangle\langle+|+a_{11--}|1\rangle\langle 1| \otimes|-\rangle\langle-| \\
& +a_{12++}|1\rangle\langle 2| \otimes|+\rangle\langle+|+a_{12+-}|1\rangle\langle 2| \otimes|+\rangle\langle-| \\
& +a_{12-+}|1\rangle\langle 2| \otimes|-\rangle\langle+|+a_{12--}|1\rangle\langle 2| \otimes|-\rangle\langle-| \\
& +a_{21++}|2\rangle\langle 1| \otimes|+\rangle\langle+|+a_{21+-}|2\rangle\langle 1| \otimes|+\rangle\langle-|  \tag{B.3}\\
& +a_{21-+}|2\rangle\langle 1| \otimes|-\rangle\langle+|+a_{21--}|2\rangle\langle 1| \otimes|-\rangle\langle-| \\
& +a_{22++}|2\rangle\langle 2| \otimes|+\rangle\langle+|+a_{22+-}|2\rangle\langle 2| \otimes|+\rangle\langle-| \\
& +a_{22-+}|2\rangle\langle 2| \otimes|-\rangle\langle+|+a_{22--}|2\rangle\langle 2| \otimes|-\rangle\langle-| .
\end{align*}
$$

We can also enumerate all the possible states, some of these can be entangled

$$
\begin{equation*}
|\psi\rangle \in h_{1+}|1\rangle \otimes|+\rangle+h_{1-}|1\rangle \otimes|-\rangle+h_{2+}|2\rangle \otimes|+\rangle+h_{2-}|2\rangle \otimes|-\rangle . \tag{B.4}
\end{equation*}
$$

Finally, we can enumerate all the possible product states

$$
\begin{align*}
|\psi\rangle & \in\left(a_{i}|i\rangle\right) \otimes\left(b_{\beta}|\beta\rangle\right) \\
& =a_{1} b_{+}|1\rangle \otimes|+\rangle+a_{1} b_{-}|1\rangle \otimes|-\rangle+a_{2} b_{+}|2\rangle \otimes|+\rangle+a_{2} b_{-}|2\rangle \otimes|-\rangle . \tag{B.5}
\end{align*}
$$

In this simpler example, we have the same dimensionality for both the sets of direct product kets and the ones formed by arbitrary superposition of the composite ket basis elements, but that does not mean that this rules out entanglement. Suppose that, as the product of some operator, we end up with a ket

$$
\begin{equation*}
|\psi\rangle=|1\rangle \otimes|+\rangle+|2\rangle \otimes|-\rangle . \tag{B.6}
\end{equation*}
$$

Does this have a product representation of the following form

$$
\begin{equation*}
|\psi\rangle=\left(a_{i}|i\rangle\right) \otimes\left(b_{\beta}|\beta\rangle\right)=a_{i} b_{\beta}|i\rangle \otimes|\beta\rangle ? \tag{B.7}
\end{equation*}
$$

For this to be true we would require

$$
\begin{align*}
& a_{1} b_{+}=1 \\
& a_{2} b_{-}=1  \tag{B.8}\\
& a_{1} b_{-}=0 \\
& a_{2} b_{+}=0 .
\end{align*}
$$

However, we can not find a solution to this set of equations. We require one of $a_{1}=0$ or $b_{-}=0$ for the third equality, but such zeros generate contradictions for one of the first pair of equations.

I was deceived by an incorrect result in Mathematica, which led me to believe that the second order energy perturbation was zero (whereas part (c) of the problem asked if it was greater or lesser than zero). I started starting writing this up to show my reasoning, but our Professor quickly provided an example after class showing how this zero must be wrong, and I did not have to show him any of this.

Setup. Recall first the one dimensional particle in a box. Within the box we have to solve

$$
\begin{equation*}
\frac{P^{2}}{2 m} \psi=E \psi \tag{C.1}
\end{equation*}
$$

and find

$$
\begin{equation*}
\psi \sim e^{\frac{i}{\hbar} \sqrt{2 m E} x} \tag{C.2}
\end{equation*}
$$

With

$$
\begin{equation*}
k=\frac{\sqrt{2 m E}}{\hbar} \tag{C.3}
\end{equation*}
$$

our general state, involving terms of each sign, takes the form

$$
\begin{equation*}
\psi=A e^{i k x}+B e^{-i k x} \tag{C.4}
\end{equation*}
$$

Inserting boundary conditions gives us

$$
\left[\begin{array}{c}
\psi(-L / 2)  \tag{C.5}\\
\psi(L / 2)
\end{array}\right]\left[\begin{array}{l}
e^{-i k \frac{L}{2}}+e^{i k \frac{L}{2}} \\
e^{i k \frac{L}{2}}+e^{-i k \frac{L}{2}}
\end{array}\right]\left[\begin{array}{l}
A \\
B
\end{array}\right]
$$

The determinant is zero

$$
\begin{equation*}
e^{-i k L}-e^{i k L}=0 \tag{C.6}
\end{equation*}
$$

which provides our constraint on $k$

$$
\begin{equation*}
e^{2 i k L}=1 \tag{C.7}
\end{equation*}
$$

We require $2 k L=2 \pi n$ for any integer $n$, or

$$
\begin{equation*}
k=\frac{\pi n}{L} . \tag{C.8}
\end{equation*}
$$

This quantizes the energy, and inverting eq. (C.3) gives us

$$
\begin{equation*}
E=\frac{1}{2 m}\left(\frac{\hbar \pi n}{L}\right)^{2} \tag{C.9}
\end{equation*}
$$

To complete the task of matching boundary value conditions we cheat and recall that the particular linear combinations that we need to match the boundary constraint of zero at $\pm L / 2$ were sums and differences yielding cosines and sines respectively. Since

$$
\begin{equation*}
\left.\sin \left(\frac{\pi n x}{L}\right)\right|_{x= \pm L / 2}= \pm \sin \left(\frac{\pi n}{2}\right) . \tag{C.10}
\end{equation*}
$$

So sines are the wave functions for $n=2,4, \ldots$ since $\sin (n \pi)=0$ for integer $n$. Similarly

$$
\begin{equation*}
\left.\cos \left(\frac{\pi n x}{L}\right)\right|_{x= \pm L / 2}=\cos \left(\frac{\pi n}{2}\right) . \tag{C.11}
\end{equation*}
$$

Cosine becomes zero at $\pi / 2,3 \pi / 2, \cdots$, so our wave function is the cosine for $n=1,3,5, \cdots$. Normalizing gives us

$$
\psi_{n}(x)=\sqrt{\frac{2}{L}} \begin{cases}\cos \left(\frac{\pi n x}{L}\right) & n=1,3,5, \cdots  \tag{C.12}\\ \sin \left(\frac{\pi n x}{L}\right) & n=2,4,6, \cdots\end{cases}
$$

Two non-interacting particles. Three lowest energy levels and degeneracies. Forming the Hamiltonian for two particles in the box without interaction, we have within the box

$$
\begin{equation*}
H=\frac{P_{1}^{2}}{2 m}+\frac{P_{2}^{2}}{2 m}, \tag{C.13}
\end{equation*}
$$

we can apply separation of variables, and it becomes clear that our wave functions have the form

$$
\begin{equation*}
\psi_{n m}\left(x_{1}, x_{2}\right)=\psi_{n}\left(x_{1}\right) \psi_{m}\left(x_{2}\right) \tag{C.14}
\end{equation*}
$$

Plugging in

$$
\begin{equation*}
H \psi=E \psi, \tag{C.15}
\end{equation*}
$$

supplies the energy levels for the two particle wavefunction, giving

$$
\begin{align*}
H \psi_{n m} & =\frac{\hbar^{2}}{2 m}\left(\left(\frac{\pi n}{L}\right)^{2}+\left(\frac{\pi m}{L}\right)^{2}\right) \psi_{n m} \\
& =\frac{1}{2 m}\left(\frac{\hbar \pi}{L}\right)^{2}\left(n^{2}+m^{2}\right) \psi_{n m} \tag{C.16}
\end{align*}
$$

Letting $n, m$ each range over $[1,3]$ for example we find
Table C.1: Varying the wave numbers.

| $n$ | $m$ | $n^{2}+m^{2}$ |
| :--- | :--- | :--- |
| 1 | 1 | 2 |
| 1 | 2 | 5 |
| 1 | 3 | 10 |
| 2 | 1 | 5 |
| 2 | 2 | 8 |
| 2 | 3 | 13 |
| 3 | 1 | 10 |
| 3 | 2 | 13 |
| 3 | 3 | 18 |

It is clear that our lowest energy levels are

$$
\begin{gather*}
\frac{1}{m}\left(\frac{\hbar \pi}{L}\right)^{2} \\
\frac{5}{2 m}\left(\frac{\hbar \pi}{L}\right)^{2}  \tag{C.17}\\
\frac{4}{m}\left(\frac{\hbar \pi}{L}\right)^{2}
\end{gather*}
$$

with degeneracies $1,2,1$ respectively.

Ground state energy with interaction perturbation to first order. With $c_{0}$ positive and an interaction potential of the form

$$
\begin{equation*}
U\left(X_{1}, X_{2}\right)=-c_{0} \delta\left(X_{1}-X_{2}\right) \tag{C.18}
\end{equation*}
$$

The second order perturbation of the ground state energy is

$$
\begin{equation*}
E=E_{11}^{(0)}+H_{11 ; 11}^{\prime}+\sum_{n m \neq 11} \frac{\left|H_{11 ; 11}^{\prime}\right|^{2}}{E_{11}-E_{n m}}, \tag{C.19}
\end{equation*}
$$

where

$$
\begin{equation*}
E_{11}^{(0)}=\frac{1}{m}\left(\frac{\hbar \pi}{L}\right)^{2} \tag{C.20}
\end{equation*}
$$

and

$$
\begin{equation*}
H_{n m ; a b}^{\prime}=-c_{0}\left\langle\psi_{n m}\right| \delta\left(X_{1}-X_{2}\right)\left|\psi_{a b}\right\rangle . \tag{C.21}
\end{equation*}
$$

To proceed, we need to expand the matrix element

$$
\begin{align*}
\left\langle\psi_{n m}\right| & \delta\left(X_{1}-X_{2}\right)\left|\psi_{a b}\right\rangle \\
& =\int d x_{1} d x_{2} d y_{1} d y_{2}\left\langle\psi_{n m} \mid x_{1} x_{2}\right\rangle\left\langle x_{1} x_{2}\right| \delta\left(X_{1}-X_{2}\right)\left|y_{1} y_{2}\right\rangle\left\langle y_{1} y_{2} \mid \psi_{a b}\right\rangle \\
& =\int d x_{1} d x_{2} d y_{1} d y_{2}\left\langle\psi_{n m} \mid x_{1} x_{2}\right\rangle \delta\left(x_{1}-x_{2}\right) \delta^{2}(\mathbf{x}-\mathbf{y})\left\langle y_{1} y_{2} \mid \psi_{a b}\right\rangle \\
& =\int d x_{1} d x_{2}\left\langle\psi_{n m} \mid x_{1} x_{2}\right\rangle \delta\left(x_{1}-x_{2}\right)\left\langle x_{1} x_{2} \mid \psi_{a b}\right\rangle \\
& =\int_{-L / 2}^{L / 2} d x \psi_{n m}(x, x) \psi_{a b}(x, x) . \tag{C.22}
\end{align*}
$$

So, for our first order calculation we need

$$
\begin{align*}
H_{11 ; 11}^{\prime} & =-c_{0} \int_{-L / 2}^{L / 2} d x \psi_{11}(x, x) \psi_{11}(x, x) \\
& =\frac{4}{L^{2}} \int_{-L / 2}^{L / 2} d x \cos ^{4}(\pi x / L)  \tag{C.23}\\
& =-\frac{3 c_{0}}{2 L} .
\end{align*}
$$

For the second order perturbation of the energy, it is clear that this will reduce the first order approximation for each matrix element that is nonzero. Attempting that calculation with Mathematica however, is deceiving, since Mathematica reports these all as zero after FullSimplify. It appears, that as used, it does not allow for $m=n$ and $m=n \pm 1$ constraints properly where the denominators of the unsimplified integrals go zero, and more careful evaluation is required. See [13] for details.

Professor Sipe's adiabatic perturbation and that of the text [4] in §17.5.1 and $\S 17.5$.2 use different notation for $\gamma_{m}$ and take a slightly different approach. We can find Prof Sipe's final result with a bit less work, if a hybrid of the two methods is used. Our starting point is the same, we have a time dependent slowly varying Hamiltonian

$$
\begin{equation*}
H=H(t) \tag{D.1}
\end{equation*}
$$

where our perturbation starts at some specific time from a given initial state

$$
\begin{equation*}
H(t)=H_{0}, \quad t \leq 0 \tag{D.2}
\end{equation*}
$$

We assume that instantaneous eigenkets can be found, satisfying

$$
\begin{equation*}
H(t)|n(t)\rangle=E_{n}(t)|n(t)\rangle \tag{D.3}
\end{equation*}
$$

Here I will use $|n\rangle \equiv|n(t)\rangle$ instead of the $\left|\hat{\psi}_{n}(t)\right\rangle$ that we used in class because its easier to write. Now suppose that we have some arbitrary state, expressed in terms of the instantaneous basis kets $|n\rangle$

$$
\begin{equation*}
|\psi\rangle=\sum_{n} \bar{b}_{n}(t) e^{-i \alpha_{n}+i \beta_{n}}|n\rangle \tag{D.4}
\end{equation*}
$$

where

$$
\begin{equation*}
\alpha_{n}(t)=\frac{1}{\hbar} \int_{0}^{t} d t^{\prime} E_{n}\left(t^{\prime}\right) \tag{D.5}
\end{equation*}
$$

Here I have used $\beta_{n}$ instead of $\gamma_{n}$ (as in the text) to avoid conflicting with the lecture notes, where this $\beta_{n}$ is a factor to be determined. For this state, we have at the time just before the perturbation

$$
\begin{equation*}
|\psi(0)\rangle=\sum_{n} \bar{b}_{n}(0) e^{-i \alpha_{n}(0)+i \beta_{n}(0)}|n(0)\rangle \tag{D.6}
\end{equation*}
$$

The question to answer is: How does this particular state evolve? Another question, for those that do not like sneaky bastard derivations, is where did
that magic factor of $e^{-i \alpha_{n}}$ come from in our superposition state? We will see after we start taking derivatives that this is what we need to cancel the $H(t)|n\rangle$ in Schrödinger's equation. Proceeding to plug into the evolution identity we have

$$
\begin{align*}
0 & =\langle m|\left(i \hbar \frac{d}{d t}-H(t)\right)|\psi\rangle \\
& =\langle m|\left(\sum_{n} e^{-i \alpha_{n}+i \beta_{n}}(i \hbar)\left(\frac{d \bar{b}_{n}}{d t}+\bar{b}_{n}\left(-i \frac{E_{n}}{\hbar}+i \dot{\beta}_{m}\right)\right)|n\rangle+i \hbar \bar{b}_{n} \frac{d}{d t}|n\rangle-E_{n} \bar{b}_{n}|n\rangle\right) \\
& =e^{-i \alpha_{m}+i \beta_{m}}(i \hbar) \frac{d \bar{b}_{m}}{d t}+e^{-i \alpha_{m}+i \beta_{m}}(i \hbar) i \dot{\beta}_{m} \bar{b}_{m}+i \hbar \sum_{n} \bar{b}_{n}\langle m| \frac{d}{d t}|n\rangle e^{-i \alpha_{n}+i \beta_{n}} \\
& \sim \frac{d \bar{b}_{m}}{d t}+i \dot{\beta}_{m} \bar{b}_{m}+\sum_{n} e^{-i \alpha_{n}+i \beta_{n}} e^{i \alpha_{m}-i \beta_{m}} \bar{b}_{n}\langle m| \frac{d}{d t}|n\rangle \\
& =\frac{d \bar{b}_{m}}{d t}+i \dot{\beta}_{m} \bar{b}_{m}+\bar{b}_{m}\langle m| \frac{d}{d t}|m\rangle+\sum_{n \neq m} e^{-i \alpha_{n}+i \beta_{n}} e^{i \alpha_{m}-i \beta_{m}} \bar{b}_{n}\langle m| \frac{d}{d t}|n\rangle . \tag{D.7}
\end{align*}
$$

We are free to pick $\beta_{m}$ to kill the second and third terms

$$
\begin{equation*}
0=i \dot{\beta}_{m} \bar{b}_{m}+\bar{b}_{m}\langle m| \frac{d}{d t}|m\rangle, \tag{D.8}
\end{equation*}
$$

or

$$
\begin{equation*}
\dot{\beta}_{m}=i\langle m| \frac{d}{d t}|m\rangle, \tag{D.9}
\end{equation*}
$$

which after integration is

$$
\begin{equation*}
\beta_{m}(t)=i \int_{0}^{t} d t^{\prime}\left\langle m\left(t^{\prime}\right)\right| \frac{d}{d t^{\prime}}|m(t)\rangle . \tag{D.10}
\end{equation*}
$$

In the lecture notes this was written as

$$
\begin{equation*}
\Gamma_{m}(t)=i\langle m(t)| \frac{d}{d t}|m(t)\rangle, \tag{D.11}
\end{equation*}
$$

so that

$$
\begin{equation*}
\beta_{m}(t)=\int_{0}^{t} d t^{\prime} \Gamma_{m}\left(t^{\prime}\right) \tag{D.12}
\end{equation*}
$$

As in class we can observe that this is a purely real function. We are left with

$$
\begin{equation*}
\frac{d \bar{b}_{m}}{d t}=-\sum_{n \neq m} \bar{b}_{n} e^{-i \alpha_{n m}+i \beta_{n m}}\langle m| \frac{d}{d t}|n\rangle, \tag{D.13}
\end{equation*}
$$

where

$$
\begin{align*}
\alpha_{n m} & =\alpha_{n}-\alpha_{m}  \tag{D.14}\\
\beta_{n m} & =\beta_{n}-\beta_{m} .
\end{align*}
$$

The task is now to find solutions for these $\bar{b}_{m}$ coefficients, and we can refer to the class notes for that without change.

Motivation. In lecture 9, Prof Sipe developed the equations governing the evolution of the coefficients of a given state for an adiabatically changing Hamiltonian. He also indicated that we could do an approximation, finding the evolution of an initially pure state in powers of $\lambda$ (like we did for the solutions of a non-time dependent perturbed Hamiltonian $H=H_{0}+\lambda H^{\prime}$ ). I tried doing that a couple of times and always ended up going in circles. I will show that here and also develop an expansion in time up to second order as an alternative, which appears to work out nicely.

Review. We assumed that an adiabatically changing Hamiltonian was known with instantaneous eigenkets governed by

$$
\begin{equation*}
H(t)\left|\hat{\psi}_{n}(t)\right\rangle=\hbar \omega_{n}\left|\hat{\psi}_{n}(t)\right\rangle \tag{E.1}
\end{equation*}
$$

The problem was to determine the time evolutions of the coefficients $\bar{b}_{n}(t)$ of some state $|\psi(t)\rangle$, and this was found to be

$$
\begin{align*}
|\psi(t)\rangle & =\sum_{n} \bar{b}_{n}(t) e^{-i \gamma_{n}(t)}\left|\hat{\psi}_{n}(t)\right\rangle \\
\gamma_{s}(t) & =\int_{0}^{t} d t^{\prime}\left(\omega_{s}\left(t^{\prime}\right)-\Gamma_{s}\left(t^{\prime}\right)\right)  \tag{E.2}\\
\Gamma_{s}(t) & =i\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{s}(t)\right\rangle
\end{align*}
$$

where the $\bar{b}_{s}(t)$ coefficient must satisfy the set of LDEs

$$
\begin{equation*}
\frac{d \bar{b}_{s}(t)}{d t}=-\sum_{n \neq s} \bar{b}_{n}(t) e^{i \gamma_{s n}(t)}\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle \tag{E.3}
\end{equation*}
$$

where

$$
\begin{equation*}
\gamma_{s n}(t)=\gamma_{s}(t)-\gamma_{n}(t) \tag{E.4}
\end{equation*}
$$

Solving these in general does not look terribly fun, but perhaps we can find an explicit solution for all the $\bar{b}_{s}$ 's, if we simplify the problem some-
what. Suppose that our initial state is found to be in the $m$ th energy level at the time before we start switching on the changing Hamiltonian.

$$
\begin{equation*}
|\psi(0)\rangle=\bar{b}_{m}(0)\left|\hat{\psi}_{m}(0)\right\rangle . \tag{E.5}
\end{equation*}
$$

We therefore require (up to a phase factor)

$$
\begin{align*}
& \bar{b}_{m}(0)=1 \\
& \bar{b}_{s}(0)=0 \quad \text { if } s \neq m . \tag{E.6}
\end{align*}
$$

Equivalently we can write

$$
\begin{equation*}
\bar{b}_{s}(0)=\delta_{m s} \tag{E.7}
\end{equation*}
$$

Going in circles with a $\lambda$ expansion. In class it was hinted that we could try a $\lambda$ expansion of the following form to determine a solution for the $\bar{b}_{s}$ coefficients at later times

$$
\begin{equation*}
\bar{b}_{s}(t)=\delta_{m s}+\lambda \bar{b}_{s}^{(1)}(t)+\cdots \tag{E.8}
\end{equation*}
$$

I was not able to figure out how to make that work. Trying this first to first order, and plugging in, we find

$$
\begin{equation*}
\lambda \frac{d}{d t} \bar{b}_{s}^{(1)}(t)=-\sum_{n \neq s}\left(\delta_{m n}+\lambda \bar{b}_{n}^{(1)}(t)\right) e^{i \gamma_{s n}(t)}\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle, \tag{E.9}
\end{equation*}
$$

equating powers of $\lambda$ yields two equations

$$
\begin{align*}
\frac{d}{d t} \bar{b}_{s}^{(1)}(t) & =-\sum_{n \neq s} \bar{b}_{n}^{(1)}(t) e^{i \gamma_{s n}(t)}\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle \\
0 & =-\sum_{n \neq s} \delta_{m n} e^{i \gamma_{s n}(t)}\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle . \tag{E.10}
\end{align*}
$$

Observe that the first identity is exactly what we started with in eq. (E.3), but has just replaced the $\bar{b}_{n}$ 's with $\bar{b}_{n}^{(1)}$ 's. Worse is that the second equation is only satisfied for $s=m$, and for $s \neq m$ we have

$$
\begin{equation*}
0=-e^{i \gamma_{s m}(t)}\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{m}(t)\right\rangle . \tag{E.11}
\end{equation*}
$$

So this $\lambda$ power series only appears to work if we somehow had $\left|\hat{\psi}_{s}(t)\right\rangle$ always orthonormal to the derivative of $\left|\hat{\psi}_{m}(t)\right\rangle$. Perhaps this could be done if the Hamiltonian was also expanded in powers of $\lambda$, but such a beastie seems foreign to the problem. Note that we do not even have any explicit dependence on the Hamiltonian in the final $\bar{b}_{n}$ differential equations, as we would probably need for such an expansion to work out.

A Taylor series expansion in time. What we can do is to expand the $\bar{b}_{n}$ 's in a power series parameterized by time. That is, again, assuming we started with energy equal to $\hbar \omega_{m}$, form

$$
\begin{equation*}
\bar{b}_{s}(t)=\delta_{s m}+\frac{t}{1!}\left(\left.\frac{d}{d t} \bar{b}_{s}(t)\right|_{t=0}\right)+\frac{t^{2}}{2!}\left(\left.\frac{d^{2}}{d t^{2}} \bar{b}_{s}(t)\right|_{t=0}\right)+\cdots \tag{E.12}
\end{equation*}
$$

The first order term we can grab right from eq. (E.3) and find

$$
\begin{align*}
\left.\frac{d \bar{b}_{s}(t)}{d t}\right|_{t=0} & =-\left.\sum_{n \neq s} \bar{b}_{n}(0)\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle\right|_{t=0} \\
& =-\left.\sum_{n \neq s} \delta_{n m}\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle\right|_{t=0}  \tag{E.13}\\
& =\left\{\begin{array}{lr}
0 & s=m \\
-\left.\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{m}(t)\right\rangle\right|_{t=0} & s \neq m
\end{array}\right.
\end{align*}
$$

Let us write

$$
\begin{align*}
|n\rangle & =\left|\hat{\psi}_{n}(0)\right\rangle \\
\left|n^{\prime}\right\rangle & =\left.\frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle\right|_{t=0} \tag{E.14}
\end{align*}
$$

So we can write

$$
\begin{equation*}
\left.\frac{d \bar{b}_{s}(t)}{d t}\right|_{t=0}=-\left(1-\delta_{s m}\right)\left\langle s \mid m^{\prime}\right\rangle \tag{E.15}
\end{equation*}
$$

and form, to first order in time our approximation for the coefficient is

$$
\begin{equation*}
\bar{b}_{s}(t)=\delta_{s m}-t\left(1-\delta_{s m}\right)\left\langle s \mid m^{\prime}\right\rangle \tag{E.16}
\end{equation*}
$$

Let us do the second order term too. For that we have

$$
\begin{align*}
& \left.\frac{d^{2}}{d t^{2}} \bar{b}_{s}(t)\right|_{t=0} \\
& =-\left.\sum_{n \neq s}\left(\left(\frac{d}{d t} \bar{b}_{n}(t)+\delta_{n m} i \frac{d \gamma_{s n}(t)}{d t}\right)\left\langle s \mid n^{\prime}\right\rangle+\delta_{n m} \frac{d}{d t}\left(\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle\right)\right)\right|_{t=0} \tag{E.17}
\end{align*}
$$

For the $\gamma_{s n}$ derivative we note that

$$
\begin{equation*}
\left.\frac{d}{d t} \gamma_{s}(t)\right|_{t=0}=\omega_{s}(0)-i\left\langle s \mid s^{\prime}\right\rangle \tag{E.18}
\end{equation*}
$$

so we have

$$
\begin{align*}
& \left.\frac{d^{2}}{d t^{2}} \bar{b}_{s}(t)\right|_{t=0} \\
& =-\sum_{n \neq s}\left(-\left(1-\delta_{n m}\right)\left\langle n \mid m^{\prime}\right\rangle+\delta_{n m} i\left(\omega_{s n}(0)-i\left\langle s \mid s^{\prime}\right\rangle+i\left\langle n \mid n^{\prime}\right\rangle\right)\right)\left\langle s \mid n^{\prime}\right\rangle \\
& \quad+\delta_{n m}\left(\left\langle s^{\prime} \mid n^{\prime}\right\rangle+\left\langle s \mid n^{\prime \prime}\right\rangle\right) \tag{E.19}
\end{align*}
$$

Again for $s=m$, all terms are killed. That is somewhat surprising, but suggests that we will need to normalize the coefficients after the perturbation calculation, since we have unity for one of them. For $s \neq m$ we have

$$
\begin{align*}
\left.\frac{d^{2}}{d t^{2}} \bar{b}_{s}(t)\right|_{t=0}= & \sum_{n \neq s}\left(\left\langle n \mid m^{\prime}\right\rangle-\delta_{n m} i\left(\omega_{s n}(0)-i\left\langle s \mid s^{\prime}\right\rangle+i\left\langle n \mid n^{\prime}\right\rangle\right)\right)\left\langle s \mid n^{\prime}\right\rangle \\
& -\delta_{n m}\left(\left\langle s^{\prime} \mid n^{\prime}\right\rangle+\left\langle s \mid n^{\prime \prime}\right\rangle\right) \\
= & -i\left(\omega_{s m}(0)-i\left\langle s \mid s^{\prime}\right\rangle+i\left\langle m \mid m^{\prime}\right\rangle\right)\left\langle s \mid m^{\prime}\right\rangle-\left(\left\langle s^{\prime} \mid m^{\prime}\right\rangle+\left\langle s \mid m^{\prime \prime}\right\rangle\right) \\
& +\sum_{n \neq s}\left\langle n \mid m^{\prime}\right\rangle\left\langle s \mid n^{\prime}\right\rangle \tag{E.20}
\end{align*}
$$

So we have, for $s \neq m$

$$
\begin{align*}
\left.\frac{d^{2}}{d t^{2}} \bar{b}_{s}(t)\right|_{t=0}= & \left(\left\langle m \mid m^{\prime}\right\rangle-\left\langle s \mid s^{\prime}\right\rangle\right)\left\langle s \mid m^{\prime}\right\rangle-i \omega_{s m}(0)\left\langle s \mid m^{\prime}\right\rangle  \tag{E.21}\\
& -\left\langle s^{\prime} \mid m^{\prime}\right\rangle-\left\langle s \mid m^{\prime \prime}\right\rangle+\sum_{n \neq s}\left\langle n \mid m^{\prime}\right\rangle\left\langle s \mid n^{\prime}\right\rangle
\end{align*}
$$

It is not particularly illuminating looking, but possible to compute, and we can use it to form a second order approximate solution for our perturbed state.

$$
\begin{align*}
& \bar{b}_{s}(t)=\delta_{s m}-t\left(1-\delta_{s m}\right)\left\langle s \mid m^{\prime}\right\rangle+\left(1-\delta_{s m}\right) \times \\
& \quad\left(\left(\left\langle m \mid m^{\prime}\right\rangle-\left\langle s \mid s^{\prime}\right\rangle\right)\left\langle s \mid m^{\prime}\right\rangle-i \omega_{s m}(0)\left\langle s \mid m^{\prime}\right\rangle\right.  \tag{E.22}\\
& \left.\quad-\left\langle s^{\prime} \mid m^{\prime}\right\rangle-\left\langle s \mid m^{\prime \prime}\right\rangle+\sum_{n \neq s}\left\langle n \mid m^{\prime}\right\rangle\left\langle s \mid n^{\prime}\right\rangle\right) \frac{t^{2}}{2}
\end{align*}
$$

New info. How to do the $\lambda$ expansion? Asking about this, Federico nicely explained. "The reason why you are going in circles when trying the lambda expansion is because you are not assuming the term $\langle\psi(t)|(d / d t)|\psi(t)\rangle$ to be of order lambda. This has to be assumed, otherwise it does not make sense at all trying a perturbative approach. This assumption means that the coupling between the level $s$ and the other levels is assumed to be small because the time dependent part of the Hamiltonian is small or changes slowly with time. Making a Taylor expansion in time would be sensible only if you are interested in a short interval of time. The lambdaexpansion approach would work for any time as long as the time dependent piece of the Hamiltonian does not change wildly or is too big." In the tutorial he outlined another way to justify this. We have written so far

$$
H= \begin{cases}H(t) & t>0  \tag{E.23}\\ H_{0} & t<0\end{cases}
$$

where $H(0)=H_{0}$. We can make this explicit, and introduce a $\lambda$ factor into the picture if we write

$$
\begin{equation*}
H(t)=H_{0}+\lambda H^{\prime}(t) \tag{E.24}
\end{equation*}
$$

where $H_{0}$ has no time dependence, so that our Hamiltonian is then just the "steady-state" system for $\lambda=0$. Now recall the method from [2] that we can use to relate our bra-derivative-ket to the Hamiltonian. Taking derivatives of the energy identity, braketed between two independent kets ( $m \neq n$ ) we have

$$
\begin{align*}
0 & =\left\langle\hat{\psi}_{m}(t)\right| \frac{d}{d t}\left(H(t)\left|\hat{\psi}_{n}(t)\right\rangle-\hbar \omega_{n}\left|\hat{\psi}_{n}(t)\right\rangle\right) \\
& =\left\langle\hat{\psi}_{m}(t)\right|\left(\frac{d H(t)}{d t}\left|\hat{\psi}_{n}(t)\right\rangle+H(t) \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle-\hbar \frac{d \omega_{n}}{d t}\left|\hat{\psi}_{n}(t)\right\rangle-\hbar \omega_{n} \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle\right) \\
& =\hbar\left(\omega_{m}-\omega_{n}\right)\left\langle\hat{\psi}_{m}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle-\hbar \frac{d \omega_{n}}{d t} \delta_{m n}+\left\langle\hat{\psi}_{m}(t)\right| \frac{d H(t)}{d t}\left|\hat{\psi}_{n}(t)\right\rangle . \tag{E.25}
\end{align*}
$$

So for $m \neq n$ we find a dependence between the bra-derivative-ket and the time derivative of the Hamiltonian

$$
\begin{equation*}
\left\langle\hat{\psi}_{m}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle=\frac{\left\langle\hat{\psi}_{m}(t)\right| \frac{d H(t)}{d t}\left|\hat{\psi}_{n}(t)\right\rangle}{\hbar\left(\omega_{n}-\omega_{m}\right)} \tag{E.26}
\end{equation*}
$$

Referring back to eq. (E.24) we see the $\lambda$ dependence in this quantity, coming directly from the $\lambda$ dependence imposed on the time dependent part of the Hamiltonian

$$
\begin{equation*}
\left\langle\hat{\psi}_{m}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle=\lambda \frac{\left\langle\hat{\psi}_{m}(t)\right| \frac{d H^{\prime}(t)}{d t}\left|\hat{\psi}_{n}(t)\right\rangle}{\hbar\left(\omega_{n}-\omega_{m}\right)} \tag{E.27}
\end{equation*}
$$

Given this $\lambda$ dependence, let us revisit the perturbation attempt of eq. (E.9). Our first order factors of $\lambda$ are now

$$
\begin{align*}
\frac{d}{d t} \bar{b}_{s}^{(1)}(t) & =-\sum_{n \neq s} \delta_{m n} e^{i \gamma_{s n}(t)}\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle \\
& = \begin{cases}0 & \text { if } m=s \\
-e^{i \gamma_{s m}(t)}\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{m}(t)\right\rangle & \text { if } m \neq s\end{cases} \tag{E.28}
\end{align*} .
$$

So we find to first order

$$
\begin{equation*}
\bar{b}_{s}(t)=\delta_{m s}(1+\lambda \text { constant })-\left(1-\delta_{m s}\right) \lambda \int_{0}^{t} d t^{\prime} e^{i \gamma_{s m}\left(t^{\prime}\right)}\left\langle\hat{\psi}_{s}\left(t^{\prime}\right)\right| \frac{d}{d t^{\prime}}\left|\hat{\psi}_{m}\left(t^{\prime}\right)\right\rangle \tag{E.29}
\end{equation*}
$$

A couple observations of this result. One is that the constant factor in the $m=s$ case makes sense. This would likely be a negative contribution since we have to decrease the probability coefficient for finding our wavefunction in the $m=s$ state after perturbation, since we are increasing the probability for finding it elsewhere by changing the Hamiltonian. Also observe that since $e^{i \gamma_{s m}} \sim 0$ for small $t$ this is consistent with the first order Taylor series expansion where we found our first order contribution was

$$
\begin{equation*}
-\left(1-\delta_{m s}\right) t\left\langle\hat{\psi}_{s}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{m}(t)\right\rangle \tag{E.30}
\end{equation*}
$$

Also note that this $-e^{i \gamma_{s m}\left(t^{\prime}\right)}\left\langle\hat{\psi}_{s}\left(t^{\prime}\right)\right| \frac{d}{d t^{\prime}}\left|\hat{\psi}_{m}\left(t^{\prime}\right)\right\rangle$ is exactly the difference from 0 that was mentioned in class when the trial solution of $\bar{b}_{s}=\delta_{s m}$ was tested by plugging it into eq. (E.3), so it is not too surprising that we should have a factor of exactly this form when we refine our approximation.

A question to consider should we wish to refine the $\lambda$ perturbation to higher than first order in $\lambda$ : is there any sort of $\lambda$ dependence in the $e^{i \gamma_{s m}}$ coming from the $\Gamma_{s m}$ term in that exponential?

## F. 1 motivation.

In class it was mentioned that to deal with perturbation around a degenerate energy eigenvalue, we needed to diagonalize the perturbing Hamiltonian. I did not follow those arguments completely, and I had like to revisit those here.

## f. 2 a four state hamiltonian.

Problem set 3, problem 1, was to calculate the energy eigenvalues for the following Hamiltonian

$$
\begin{align*}
H & =H_{0}+\lambda H^{\prime} \\
H_{0} & =\left[\begin{array}{llll}
a & 0 & 0 & 0 \\
0 & b & 0 & 0 \\
0 & 0 & c & 0 \\
0 & 0 & 0 & c
\end{array}\right]  \tag{F.1}\\
H^{\prime} & =\left[\begin{array}{llll}
\alpha & 0 & v & \eta \\
0 & \beta & 0 & \mu \\
v^{*} & 0 & \gamma & 0 \\
\eta^{*} & \mu^{*} & 0 & \delta
\end{array}\right] .
\end{align*}
$$

This is more complicated that the two state problem that are solved exactly in $\S 13.1 .1$ in the text [4], but differs from the (possibly) infinite dimensional problem that was covered in class. Unfortunately, the solution provided to this problem did not provide the illumination I expected, so let us do it again, calculating the perturbed energy eigenvalues for the degenerate levels, from scratch. Can we follow the approach used in the text for the two (only) state problem. For the two state problem, it was assumed that the perturbed solution could be expressed as a superposition
of the two states that formed the basis for the unperturbed Hilbert space. That is

$$
\begin{equation*}
|\psi\rangle=m|1\rangle+n|2\rangle . \tag{F.2}
\end{equation*}
$$

For the two state problem, assuming that the perturbed energy eigenvalue is $E$, and the unperturbed energy eigenvalue is $E^{0}$ we find

$$
\left.\begin{array}{rl}
0 & =(H-E)|\psi\rangle \\
& =\left(H_{0}+\lambda H^{\prime}\right)|\psi\rangle-E|\psi\rangle \\
& =\left(H_{0}+\lambda H^{\prime}\right)(m|1\rangle+n|2\rangle)-E(m|1\rangle+n|2\rangle) \\
& =\lambda H^{\prime}(m|1\rangle+n|2\rangle)-E^{0}(m|1\rangle+n|2\rangle)  \tag{F.3}\\
& =\left(-E^{0}+\lambda H^{\prime}\right)[|1\rangle \quad|2\rangle
\end{array}\right]\left[\begin{array}{l}
m \\
n
\end{array}\right] .
$$

Left multiplying by the brakets we find

$$
\begin{align*}
0 & =\left[\begin{array}{c}
\langle 1| \\
\langle 2|
\end{array}\right](H-E)|\psi\rangle \\
& =\left(\left(E^{0}-E\right) I+\lambda\left[\begin{array}{ll}
\langle 1| H^{\prime}|1\rangle & \langle 1| H^{\prime}|2\rangle \\
\langle 2| H^{\prime}|1\rangle & \langle 2| H^{\prime}|2\rangle
\end{array}\right]\right)\left[\begin{array}{l}
m \\
n
\end{array}\right], \tag{F.4}
\end{align*}
$$

or

$$
\left(\left(E^{0}-E\right) I+\lambda\left[H_{i j}^{\prime}\right]\right)\left[\begin{array}{l}
m  \tag{F.5}\\
n
\end{array}\right]=0
$$

Observe that there was no assumption about the dimensionality of $H_{0}$ and $H^{\prime}$ here, just that the two degenerate energy levels had eigenvalues $E^{0}$ and a pair of eigenkets $|1\rangle$ and $|2\rangle$ such that $H_{0}|i\rangle=E^{0}|i\rangle, i \in[1,2]$. It is clear that we can use a similar argument for any degeneracy degree. It is also clear how to proceed, since we have what almost amounts to a characteristic equation for the degenerate subspace of Hilbert space for the problem. Because $H^{\prime}$ is Hermitian, a diagonalization

$$
\begin{align*}
H^{\prime} & =U^{*} D U \\
D & =\left[\mathcal{H}^{\prime}{ }_{i} \delta_{i j}\right], \tag{F.6}
\end{align*}
$$

can be found. To solve for $E$ we can take the determinant of the matrix factor of eq. (F.5), and because $I=U^{*} U$ we have

$$
\left.\begin{align*}
0 & =\mid\left(E^{0}-E\right) U^{*} I U+\lambda U^{*} D U
\end{align*} \right\rvert\,
$$

So our energy eigenvalues associated with the perturbed state are (exactly)

$$
\begin{equation*}
E=E^{0}+\lambda \mathcal{H}^{\prime}{ }_{1}, E^{0}+\lambda \mathcal{H}^{\prime}{ }_{2} . \tag{F.8}
\end{equation*}
$$

It is a bit curious seeming that only the energy eigenvalues associated with the degeneracy play any part in this result, but there is some intuitive comfort in this idea. Without the perturbation, we can not do an energy measurement that would distinguish one or the other of the eigenkets for the degenerate energy level, so it does not seem unreasonable that a perturbed energy level close to the original can be formed by superposition of these two states, and thus the perturbed energy eigenvalue for the new system would then be related to only those degenerate levels.

Observe that in the problem set three problem we had a diagonal initial Hamiltonian $H_{0}$, that does not have an impact on the argument above, since that portion of the Hamiltonian only has a diagonal contribution to the result found in eq. (F.5), since the identity $H_{0}|i\rangle=c|i\rangle, i \in[3,4]$ removes any requirement to know the specifics of that portion of the matrix element of $H_{0}$.

## F. 3 GENERALIZING SLIGHTLY.

Let us work with a system that has kets using an explicit degeneracy index

$$
\begin{equation*}
H_{0}\left|m \alpha_{m}\right\rangle=E_{m}^{0}\left|m \alpha_{m}\right\rangle, \quad \alpha_{m}=1, \cdots, \gamma_{m}, m \in[1, N] . \tag{F.9}
\end{equation*}
$$

## Example:

$$
\left|m \alpha_{m}\right\rangle \in|11\rangle
$$

$|41\rangle,|42\rangle,|43\rangle$.
Again we seek to find the energy eigenvalues of the new system

$$
\begin{equation*}
H=H_{0}+\lambda H^{\prime} . \tag{F.11}
\end{equation*}
$$

For any $m$ with associated with a degeneracy $\left(\gamma_{m}>1\right)$ we can calculate the subspace diagonalization

$$
\begin{equation*}
\left[\langle m i| H^{\prime}|m j\rangle\right]=U_{m} D_{m} U_{m}^{\dagger}, \tag{F.12}
\end{equation*}
$$

where

$$
\begin{equation*}
U_{m} U_{m}^{\dagger}=1, \tag{F.13}
\end{equation*}
$$

and $D_{m}$ is diagonal

$$
\begin{equation*}
D_{m}=\left[\delta_{i j} \mathcal{H}_{m, i}^{\prime}\right] \tag{F.14}
\end{equation*}
$$

This is not a diagonalizing transformation in the usual sense. Putting it together into block matrix form, we can write

$$
U=\left[\begin{array}{llll}
U_{1} & & &  \tag{F.15}\\
& U_{2} & & \\
& & \ddots & \\
& & & U_{N}
\end{array}\right],
$$

and find that a similarity transformation using this change of basis matrix puts all the block matrices along the diagonal into diagonal form, but leaves the rest possibly non-zero

$$
U^{\dagger}\left[\left\langle m \alpha_{m_{i}}\right| H^{\prime}|m j\rangle\left|m \alpha_{m_{j}}\right\rangle,\right] U=\left[\begin{array}{cccc}
D_{1} & x & x & x  \tag{F.16}\\
x & D_{2} & x & x \\
x & x & \ddots & x \\
x & x & x & D_{N}
\end{array}\right]
$$

A five level system with two pairs of degenerate levels. Let us do this explicitly using a specific degeneracy example, supposing that we have a non-degenerate ground state, and two pairs doubly degenerate next energy levels. That is

$$
\begin{align*}
\left|m \alpha_{m}\right\rangle \in & |11\rangle \\
& |21\rangle,|22\rangle  \tag{F.17}\\
& |31\rangle,|32\rangle .
\end{align*}
$$

Our change of basis matrix is

$$
U=\left[\begin{array}{c:cc:cc}
1 & 0 & 0 & 0 & 0  \tag{F.18}\\
\hdashline 0 & U_{2} & 0 & 0 \\
0 & : & 0 & 0 \\
\hdashline 0 & 0 & 0 & U_{3} \\
0 & 0 & 0 &
\end{array}\right] .
$$

We would like to calculate

$$
\begin{equation*}
U^{\dagger} H^{\prime} U \tag{F.19}
\end{equation*}
$$

Let us write this putting row and column range subscripts on our matrices to explicitly block them into multiplication compatible sized pieces

$$
\begin{align*}
U & =\left[\begin{array}{lll}
I_{11,11} & 0_{11,23} & 0_{11,45} \\
0_{23,11} & U_{23,23} & 0_{23,45} \\
0_{45,11} & 0_{45,23} & U_{45,45}
\end{array}\right] \\
H^{\prime} & =\left[\begin{array}{lll}
{H^{\prime}}_{11,11} & {H^{\prime}}_{11,23} & {H^{\prime}}_{11,45} \\
{H^{\prime}}_{23,11} & {H^{\prime}}_{23,23} & {H^{\prime}}_{23,45} \\
H^{\prime}{ }_{45,11} & {H^{\prime}}_{45,23} & {H^{\prime}}_{45,45}
\end{array}\right] . \tag{F.20}
\end{align*}
$$

The change of basis calculation then becomes

$$
\begin{align*}
& U^{\dagger} H^{\prime} U=\left[\begin{array}{ccc}
I_{11,11} & 0_{11,23} & 0_{11,45} \\
0_{23,11} & U_{23,23}^{\dagger} & 0_{23,45} \\
0_{45,11} & 0_{45,23} & U_{45,45}^{\dagger}
\end{array}\right]\left[\begin{array}{ccc}
H^{\prime}{ }_{11,11} & H^{\prime}{ }_{11,23} & H^{\prime}{ }_{11,45} \\
H^{\prime}{ }_{23,11} & H^{\prime}{ }_{23,23} & H^{\prime}{ }_{23,45} \\
H^{\prime}{ }_{45,11} & H^{\prime}{ }_{45,23} & H^{\prime}{ }_{45,45}
\end{array}\right] \times \\
& {\left[\begin{array}{lll}
I_{11,11} & 0_{11,23} & 0_{11,45} \\
0_{23,11} & U_{23,23} & 0_{23,45} \\
0_{45,11} & 0_{45,23} & U_{45,45}
\end{array}\right]} \\
& =\left[\begin{array}{ccc}
I_{11,11} & 0_{11,23} & 0_{11,45} \\
0_{23,11} & U_{23,23}^{\dagger} & 0_{23,45} \\
0_{45,11} & 0_{45,23} & U_{45,45}^{\dagger}
\end{array}\right]\left[\begin{array}{lll}
H^{\prime}{ }_{11,11} & H^{\prime}{ }_{11,23} U_{23,23} & H^{\prime}{ }_{11,45} U_{45,45} \\
H^{\prime}{ }_{23,11} & H^{\prime}{ }_{23,23} U_{23,23} & H^{\prime}{ }_{23,45} U_{45,45} \\
H^{\prime}{ }_{45,11} & H^{\prime}{ }_{45,23} U_{23,23} & H^{\prime}{ }_{45,45} U_{45,45}
\end{array}\right] \\
& =\left[\begin{array}{ccc}
H^{\prime}{ }_{11,11} & H^{\prime}{ }_{11,23} U_{23,23} & H^{\prime}{ }_{11,45} U_{45,45} \\
U_{23,23}^{\dagger} H^{\prime}{ }_{23,11} & U_{23,23}^{\dagger} H^{\prime}{ }_{23,23} U_{23,23} \\
U_{45,45}^{\dagger} H^{\prime}{ }_{45,11} & U_{45,45}^{\dagger} H^{\prime}{ }_{45,23} U_{23,23} & \left.\begin{array}{l}
U_{23,23} H^{\prime}{ }_{23,45} U_{45,45} \\
U_{45,45}^{\dagger} H^{\prime}{ }_{45,45} U_{45,45}
\end{array}\right] . ~ . ~ . ~ . ~ . ~
\end{array}\right. \tag{F.21}
\end{align*}
$$

We see that we end up with explicitly diagonal matrices along the diagonal blocks, but products that are otherwise everywhere else. In the new basis our kets become

$$
\begin{equation*}
\left|m \alpha_{m}^{\prime}\right\rangle=U^{\dagger}\left|m \alpha_{m}\right\rangle \tag{F.22}
\end{equation*}
$$

Suppose we calculate this change of basis representation for $|21\rangle$ (we have implicitly assumed above that our original basis had the ordering $\{|11\rangle|21\rangle,|22\rangle,|31\rangle,|32\rangle\})$. We find

$$
\begin{align*}
\left|21^{\prime}\right\rangle & =U^{*}|21\rangle \\
& =\left[\begin{array}{ccc}
1 & 0 & 0 \\
0 & U_{2}^{*} & 0 \\
0 & 0 & U_{3}^{*}
\end{array}\right]\left[\begin{array}{c}
0 \\
1 \\
0 \\
\hdashline- \\
0
\end{array}\right] . \tag{F.23}
\end{align*}
$$

With

$$
\begin{align*}
& U_{2}=\left[\begin{array}{ll}
U_{2,11} & U_{2,12} \\
U_{2,21} & U_{2,22}
\end{array}\right] \\
& U_{2}^{\dagger}=\left[\begin{array}{ll}
U_{2,11}^{*} & U_{2,21}^{*} \\
U_{2,12}^{*} & U_{2,22}^{*}
\end{array}\right] . \tag{F.24}
\end{align*}
$$

We find

$$
\begin{align*}
\left|21^{\prime}\right\rangle & =U^{*}\left|21^{\prime}\right\rangle \\
& =\left[\begin{array}{c}
0 \\
U_{2,11}^{*} \\
U_{2,12}^{*} \\
0 \\
0
\end{array}\right]=U_{2,11}^{*}|21\rangle+U_{2,12}^{*}|22\rangle . \tag{F.25}
\end{align*}
$$

Energy eigenvalues of the unperturbed Hamiltonian in the new basis. Generalizing this, it is clear that for a given degeneracy level, the transformed kets in the new basis are superposition of only the kets associated with that degenerate level (and the kets for the non-degenerate levels are left as is). Even better, we have for all $\left|m \alpha_{m}^{\prime}\right\rangle=U^{\dagger}\left|m \alpha_{m}\right\rangle$ that $\left|m \alpha_{m}^{\prime}\right\rangle$ remain eigenkets of the unperturbed Hamiltonian. We see that by computing the matrix element of our Hamiltonian in the full basis. Writing

$$
\begin{equation*}
F=U^{\dagger} H^{\prime} U, \tag{F.26}
\end{equation*}
$$

or

$$
\begin{equation*}
H^{\prime}=U F U^{\dagger}, \tag{F.27}
\end{equation*}
$$

where $F$ has been shown to have diagonal block diagonals, we can write

$$
\begin{align*}
H & =H_{0}+\lambda U F U^{\dagger} \\
& =U U^{\dagger} H_{0} U U^{\dagger}+\lambda U F U^{\dagger}  \tag{F.28}\\
& =U\left(U^{\dagger} H_{0} U+\lambda F\right) U^{\dagger} .
\end{align*}
$$

In the $\left|m \alpha_{m}^{\prime}\right\rangle$ basis, our Hamiltonian's matrix element is

$$
\begin{equation*}
H \rightarrow U^{\dagger} H_{0} U+\lambda F \tag{F.29}
\end{equation*}
$$

When $\lambda=0$, application of this Hamiltonian to the new basis kets gives

$$
\begin{align*}
H_{0}\left|m \alpha^{\prime}\right\rangle & =U^{\dagger} H_{0} U U^{\dagger}|m \alpha\rangle \\
& =U^{\dagger} H_{0}|m \alpha\rangle \\
& =U^{\dagger} \mathcal{H}_{m}^{0}|m \alpha\rangle  \tag{F.30}\\
& =\mathcal{H}_{m}^{0}\left(U^{\dagger}|m \alpha\rangle\right)
\end{align*}
$$

This is just

$$
\begin{equation*}
H_{0}\left|m \alpha^{\prime}\right\rangle=\mathcal{H}_{m}^{0}\left|m \alpha^{\prime}\right\rangle \tag{F.31}
\end{equation*}
$$

a statement that the $\left|m \alpha^{\prime}\right\rangle$ are still the energy eigenkets for the unperturbed system. This matches our expectations since we have seen that these differ from the original basis elements only for degenerate energy levels, and that these new basis elements are superpositions of only the kets for their respective degeneracy levels.

## REVIEW OF APPROXIMATION RESULTS.

G. 1 motivation.

Here I will summarize what I had put on a cheat sheet for the tests or exam, if one would be allowed. While I can derive these results, memorization unfortunately appears required for good test performance in this class, and this will give me a good reference of what to memorize.

This set of review notes covers all the approximation methods we covered except for Fermi's golden rule.

## G. 2 VARIATIONAL METHOD.

We can find an estimate of our ground state energy using

$$
\begin{equation*}
\frac{\langle\Psi| H|\Psi\rangle}{\langle\Psi \mid \Psi\rangle} \geq E_{0} \tag{G.1}
\end{equation*}
$$

## G. 3 TIME INDEPENDENT PERTURBATION.

Given a perturbed Hamiltonian and an associated solution for the unperturbed state

$$
\begin{align*}
H & =H_{0}+\lambda H^{\prime}, \quad \lambda \in[0,1] \\
H_{0}\left|\psi_{m \alpha}{ }^{(0)}\right\rangle & =E_{m}{ }^{(0)}\left|\psi_{m \alpha}{ }^{(0)}\right\rangle, \tag{G.2}
\end{align*}
$$

we assume a power series solution for the energy

$$
\begin{equation*}
E_{m}=E_{m}^{(0)}+\lambda E_{m}^{(1)}+\lambda^{2} E_{m}^{(2)}+\cdots \tag{G.3}
\end{equation*}
$$

For a non-degenerate state $\left|\psi_{m}\right\rangle=\left|\psi_{m 1}\right\rangle$, with an unperturbed value of $\left|\psi_{m}^{(0)}\right\rangle=\left|\psi_{m 1}^{(0)}\right\rangle$, we seek a power series expansion of this ket in the perturbed system

$$
\begin{align*}
\left|\psi_{m}\right\rangle= & \sum_{n, \alpha} c_{n \alpha ; m}{ }^{(0)}\left|\psi_{n \alpha}{ }^{(0)}\right\rangle+\lambda \sum_{n, \alpha} c_{n \alpha ; m}{ }^{(1)}\left|\psi_{n \alpha}{ }^{(0)}\right\rangle \\
& +\lambda^{2} \sum_{n, \alpha} c_{n \alpha ; m}{ }^{(2)}\left|\psi_{n \alpha}{ }^{(0)}\right\rangle+\cdots \\
\propto & \left|\psi_{m}{ }^{(0)}\right\rangle+\lambda \sum_{n \neq m, \alpha} \bar{c}_{n \alpha ; m}{ }^{(1)}\left|\psi_{n \alpha}{ }^{(0)}\right\rangle+\lambda^{2} \sum_{n \neq m, \alpha} \bar{c}_{n \alpha ; m}{ }^{(2)}\left|\psi_{n \alpha}{ }^{(0)}\right\rangle+\cdots \tag{G.4}
\end{align*}
$$

Any states $n \neq m$ are allowed to have degeneracy. For this case, we found to second order in energy and first order in the kets

$$
\begin{align*}
E_{m} & =E_{m}^{(0)}+\lambda H_{m 1 ; m 1}{ }^{\prime}+\lambda^{2} \sum_{n \neq m, \alpha} \frac{\left|H_{n \alpha ; m 1^{\prime}}\right|^{2}}{E_{m}^{(0)}-E_{n}^{(0)}}+\cdots \\
\left|\psi_{m}\right\rangle & \propto\left|\psi_{m}^{(0)}\right\rangle+\lambda \sum_{n \neq m, \alpha} \frac{H_{n \alpha ; m 1^{\prime}}^{\prime}}{E_{m}^{(0)}-E_{n}^{(0)}}\left|\psi_{n \alpha}^{(0)}\right\rangle+\cdots  \tag{G.5}\\
H_{n \alpha ; s \beta}^{\prime} & =\left\langle\psi_{n \alpha}^{(0)}\right| H^{\prime}\left|\psi_{s \beta}^{(0)}\right\rangle .
\end{align*}
$$

## G. 4 DEGENERACY.

When the initial energy eigenvalue $E_{m}$ has a degeneracy $\gamma_{m}>1$ we use a different approach to compute the perturbed energy eigenkets and perturbed energy eigenvalues. Writing the kets as $|m \alpha\rangle$, then we assume that the perturbed ket is a superposition of the kets in the degenerate energy level

$$
\begin{equation*}
|m \alpha\rangle^{\prime}=\sum_{i} c_{i}|m i\rangle \tag{G.6}
\end{equation*}
$$

We find that we must have

$$
\left(\left(E^{0}-E\right) I+\lambda\left[H_{m i ; m j}^{\prime}\right]\right)\left[\begin{array}{c}
c_{1}  \tag{G.7}\\
c_{2} \\
\vdots \\
c_{\gamma_{m}}
\end{array}\right]=0
$$

Diagonalizing this matrix $\left[H_{m i ; m j}^{\prime}\right]$ (a subset of the complete $H^{\prime}$ matrix element)

$$
\begin{equation*}
\left[\langle m i| H^{\prime}|m j\rangle\right]=U_{m}\left[\delta_{i j} \mathcal{H}_{m, i}^{\prime}\right] U_{m}^{\dagger}, \tag{G.8}
\end{equation*}
$$

we find, by taking the determinant, that the perturbed energy eigenvalues are in the set

$$
\begin{equation*}
E=E_{m}^{0}+\lambda \mathcal{H}_{m, i}^{\prime}, \quad i \in\left[1, \gamma_{m}\right] . \tag{G.9}
\end{equation*}
$$

To compute the perturbed kets we must work in a basis for which the block diagonal matrix elements are diagonal for all $m$, as in

$$
\begin{equation*}
\left[\langle m i| H^{\prime}|m j\rangle\right]=\left[\delta_{i j} \mathcal{H}_{m, i}^{\prime}\right] . \tag{G.10}
\end{equation*}
$$

If that is not the case, then the unitary matrices of eq. (G.8) can be computed, and the matrix

$$
U=\left[\begin{array}{llll}
U_{1} & & &  \tag{G.11}\\
& U_{2} & & \\
& & \ddots & \\
& & & U_{N}
\end{array}\right],
$$

can be formed. The kets

$$
\begin{equation*}
|\overline{m \alpha}\rangle=U^{\dagger}|m \alpha\rangle, \tag{G.12}
\end{equation*}
$$

will still be energy eigenkets of the unperturbed Hamiltonian

$$
\begin{equation*}
H_{0}|\overline{m \alpha}\rangle=E_{m}^{0}|\overline{m \alpha}\rangle, \tag{G.13}
\end{equation*}
$$

but also ensure that the partial diagonalization condition of eq. (G.8) is satisfied. In this basis, dropping overbars, the first order perturbation results found previously for perturbation about a non-degenerate state also hold, allowing us to write

$$
\begin{equation*}
|s \alpha\rangle^{\prime}=|s \alpha\rangle+\lambda \sum_{m \neq s, \beta} \frac{H^{\prime}{ }_{m \beta ; s \alpha}}{E_{s}^{(0)}-E_{m}^{(0)}}|m \beta\rangle+\cdots \tag{G.14}
\end{equation*}
$$

## G. 5 INTERACTION PICTURE.

We split of the Hamiltonian into time independent and time dependent parts, and also factorize the time evolution operator

$$
\begin{align*}
H & =H_{0}+H_{I}(t) \\
\left|\alpha_{S}\right\rangle & =e^{-i H_{0} t / \hbar}\left|\alpha_{I}(t)\right\rangle=e^{-i H_{0} t / \hbar} U_{I}(t)\left|\alpha_{I}(0)\right\rangle . \tag{G.15}
\end{align*}
$$

Plugging into Schrödinger's equation we find

$$
\begin{align*}
i \hbar \frac{d}{d t}\left|\alpha_{I}(t)\right\rangle & =H_{I}(t)\left|\alpha_{I}(t)\right\rangle \\
i \hbar \frac{d U_{I}}{d t} & =H_{I}^{\prime} U_{I}  \tag{G.16}\\
H_{I}^{\prime}(t) & =e^{i H_{0} t / \hbar} H_{I}(t) e^{-i H_{0} t / \hbar} .
\end{align*}
$$

## g. 6 time dependent perturbation.

We moved on to time dependent perturbations of the form

$$
\begin{align*}
H(t) & =H_{0}+H^{\prime}(t) \\
H_{0}\left|\psi_{n}^{(0)}\right\rangle & =\hbar \omega_{n}\left|\psi_{n}^{(0)}\right\rangle \tag{G.17}
\end{align*}
$$

where $\hbar \omega_{n}$ are the energy eigenvalues, and $\left|\psi_{n}^{(0)}\right\rangle$ the energy eigenstates of the unperturbed Hamiltonian. Use of the interaction picture led quickly to the problem of seeking the coefficients describing the perturbed state

$$
\begin{equation*}
|\psi(t)\rangle=\sum_{n} c_{n}(t) e^{-i \omega_{n} t}\left|\psi_{n}^{(0)}\right\rangle \tag{G.18}
\end{equation*}
$$

and plugging in we found

$$
\begin{align*}
i \hbar \dot{c}_{s} & =\sum_{n} H_{s n}^{\prime}(t) e^{i \omega_{s n} t} c_{n}(t) \\
\omega_{s n} & =\omega_{s}-\omega_{n}  \tag{G.19}\\
H_{s n}^{\prime}(t) & =\left\langle\psi_{s}^{(0)}\right| H^{\prime}(t)\left|\psi_{n}^{(0)}\right\rangle
\end{align*}
$$

Perturbation expansion in series. Introducing a $\lambda$ parameterized dependence in the perturbation above, and assuming a power series expansion of our coefficients

$$
\begin{align*}
H^{\prime}(t) & \rightarrow \lambda H^{\prime}(t) \\
c_{s}(t) & =c_{s}^{(0)}(t)+\lambda c_{s}^{(1)}(t)+\lambda^{2} c_{s}^{(2)}(t)+\cdots \tag{G.20}
\end{align*}
$$

we found, after equating powers of $\lambda$ a set of coupled differential equations

$$
\begin{aligned}
i \hbar \dot{c}_{s}^{(0)}(t) & =0 \\
i \hbar \dot{c}_{s}^{(1)}(t) & =\sum_{n} H_{s n}^{\prime}(t) e^{i \omega_{s n} t} c_{n}^{(0)}(t) \\
i \hbar \dot{c}_{s}^{(2)}(t) & =\sum_{n} H_{s n}^{\prime}(t) e^{i \omega_{s n} t} c_{n}^{(1)}(t) \\
& \vdots
\end{aligned}
$$

Of particular value was the expansion, assuming that we started with an initial state in energy level $m$ before the perturbation was "turned on" (ie: $\lambda=0$ )

$$
\begin{equation*}
|\psi(t)\rangle=e^{-i \omega_{m} t}\left|\psi_{m}^{(0)}\right\rangle, \tag{G.22}
\end{equation*}
$$

so that $c_{n}^{(0)}(t)=\delta_{n m}$. We then found a first order approximation for the transition probability coefficient of

$$
\begin{equation*}
i \hbar \dot{c}_{m}^{(1)}=H_{m s}^{\prime}(t) e^{i \omega_{m s} t} . \tag{G.23}
\end{equation*}
$$

## g. 7 Sudden perturbations.

The idea here is that we integrate Schrödinger's equation over the small interval containing the changing Hamiltonian

$$
\begin{equation*}
|\psi(t)\rangle=\left|\psi\left(t_{0}\right)\right\rangle+\frac{1}{i \hbar} \int_{t_{0}}^{t} H\left(t^{\prime}\right)\left|\psi\left(t^{\prime}\right)\right\rangle d t^{\prime}, \tag{G.24}
\end{equation*}
$$

and find

$$
\begin{equation*}
\left|\psi_{\text {after }}\right\rangle=\left|\psi_{\text {before }}\right\rangle . \tag{G.25}
\end{equation*}
$$

An implication is that, say, we start with a system measured in a given energy, that same system after the change to the Hamiltonian will then be in a state that is now a superposition of eigenkets from the new Hamiltonian.

## G. 8 ADIABATIC PERTURBATIONS.

Given a Hamiltonian that turns on slowly at $t=0$, a set of instantaneous eigenkets for the duration of the time dependent interval, and a representation in terms of the instantaneous eigenkets

$$
\begin{align*}
H(t) & =H_{0}, \quad t \leq 0 \\
H(t)\left|\hat{\psi}_{n}(t)\right\rangle & =E_{n}(t)\left|\hat{\psi}_{n}(t)\right\rangle \\
|\psi\rangle & =\sum_{n} \bar{b}_{n}(t) e^{-i \alpha_{n}+i \beta_{n}}\left|\hat{\psi}_{n}\right\rangle  \tag{G.26}\\
\alpha_{n}(t) & =\frac{1}{\hbar} \int_{0}^{t} d t^{\prime} E_{n}\left(t^{\prime}\right) .
\end{align*}
$$

Plugging into Schrödinger's equation we find

$$
\begin{align*}
\frac{d \bar{b}_{m}}{d t} & =-\sum_{n \neq m} \bar{b}_{n} e^{-i \gamma_{n m}}\left\langle\hat{\psi}_{m}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle \\
\gamma_{n m}(t) & =\alpha_{n}(t)-\alpha_{m}(t)-\left(\beta_{n}(t)-\beta_{m}(t)\right) \\
\beta_{n}(t) & =\int_{0}^{t} d t^{\prime} \Gamma_{n}\left(t^{\prime}\right)  \tag{G.27}\\
\Gamma_{n}(t) & =i\left\langle\hat{\psi}_{n}(t)\right| \frac{d}{d t}\left|\hat{\psi}_{n}(t)\right\rangle .
\end{align*}
$$

Here $\Gamma_{n}(t)$ is called the Berry phase.
Evolution of a given state. Given a system initially measured with energy $E_{m}(0)$ before the time dependence is "turned on"

$$
\begin{equation*}
|\psi(0)\rangle=\left|\hat{\psi}_{m}(0)\right\rangle \tag{G.28}
\end{equation*}
$$

we find that the first order Taylor series expansion for the transition probability coefficients are

$$
\begin{equation*}
\bar{b}_{s}(t)=\delta_{s m}-\left.t\left(1-\delta_{s m}\right)\left\langle\hat{\psi}_{s}(0)\right| \frac{d}{d t}\left|\hat{\psi}_{m}(t)\right\rangle\right|_{t=0} \tag{G.29}
\end{equation*}
$$

If we introduce a $\lambda$ perturbation, separating all the (slowly changing) time dependent part of the Hamiltonian $H^{\prime}$ from the non time dependent parts $H_{0}$ as in

$$
\begin{equation*}
H(t)=H_{0}+\lambda H^{\prime}(t), \tag{G.30}
\end{equation*}
$$

then we find our perturbed coefficients are

$$
\begin{align*}
\bar{b}_{s}(t)= & \delta_{m s}(1+\lambda \text { constant }) \\
& -\left(1-\delta_{m s}\right) \lambda \int_{0}^{t} d t^{\prime} e^{i \gamma_{s m}\left(t^{\prime}\right)}\left\langle\hat{\psi}_{s}\left(t^{\prime}\right)\right| \frac{d}{d t^{\prime}}\left|\hat{\psi}_{m}\left(t^{\prime}\right)\right\rangle . \tag{G.31}
\end{align*}
$$

## G. 9 wкв.

We write Schrödinger's equation as

$$
\begin{align*}
0 & =\frac{d^{2} U}{d x^{2}}+k^{2} U  \tag{G.32}\\
k^{2} & =-\kappa^{2}=\frac{2 m(E-V)}{\hbar}
\end{align*}
$$

and seek solutions of the form $U \propto e^{i \phi}$. Schrödinger's equation takes the form

$$
\begin{equation*}
-\left(\phi^{\prime}(x)\right)^{2}+i \phi^{\prime \prime}(x)+k^{2}(x)=0 \tag{G.33}
\end{equation*}
$$

Initially setting $\phi^{\prime \prime}=0$ we refine our approximation to find

$$
\begin{equation*}
\phi^{\prime}(x)=k(x) \sqrt{1+i \frac{k^{\prime}(x)}{k^{2}(x)}} \tag{G.34}
\end{equation*}
$$

To first order, this gives us

$$
\begin{equation*}
U(x) \propto \frac{1}{\sqrt{k(x)}} e^{ \pm i \int d x k(x)} \tag{G.35}
\end{equation*}
$$

What we did not cover in class, but required in the problems was the BohrSommerfeld condition described in §24.1.2 of the text [4].

$$
\begin{equation*}
\int_{x_{1}}^{x_{2}} d x \sqrt{2 m(E-V(x))}=\left(n+\frac{1}{2}\right) \pi \hbar \tag{G.36}
\end{equation*}
$$

This was found from the WKB connection formulas, themselves found my some Bessel function arguments that I have to admit that I did not understand.

## CLEBSH-GORDAN ZERO COEFFICIENTS.

h. 1 motivation.

In $\S 28.2$ of the text [4] is a statement that the Clebsh-Gordan coefficient

$$
\begin{equation*}
\left\langle m_{1} m_{2} \mid j m\right\rangle \tag{H.1}
\end{equation*}
$$

unless $m=m_{1}+m_{2}$. It appeared that it was related to the operation of $J_{z}$, but how exactly was not obvious to me. In tutorial today we hashed through this. Here is the details lying behind this statement

## h. 2 RECAP ON NOTATION.

We are taking an arbitrary two particle ket and decomposing it utilizing an insertion of a complete set of states

$$
\begin{equation*}
|j m\rangle=\sum_{m_{1}^{\prime} m_{2}^{\prime}}\left(\left|j_{1} m_{1}^{\prime}\right\rangle\left|j_{2} m_{2}^{\prime}\right\rangle\left\langle j_{1} m_{1}^{\prime}\right|\left\langle j_{2} m_{2}^{\prime}\right|\right)|j m\rangle \tag{H.2}
\end{equation*}
$$

with $j_{1}$ and $j_{2}$ fixed, this is written with the shorthand

$$
\begin{align*}
\left|j_{1} m_{1}\right\rangle\left|j_{2} m_{2}\right\rangle & =\left|m_{1} m_{2}\right\rangle  \tag{H.3}\\
\left\langle j_{1} m_{1}\right|\left\langle j_{2} m_{2}\right||j m\rangle & =\left\langle m_{1} m_{2} \mid j m\right\rangle
\end{align*}
$$

so that we write

$$
\begin{equation*}
|j m\rangle=\sum_{m_{1}^{\prime} m_{2}^{\prime}}\left|m_{1}^{\prime} m_{2}^{\prime}\right\rangle\left\langle m_{1}^{\prime} m_{2}^{\prime} \mid j m\right\rangle \tag{H.4}
\end{equation*}
$$

h. 3 THE $J_{z}$ ACTION.

We have two ways that we can apply the operator $J_{z}$ to $|j m\rangle$. One is using the sum above, for which we find

$$
\begin{align*}
J_{z}|j m\rangle & =\sum_{m_{1}^{\prime} m_{2}^{\prime}} J_{z}\left|m_{1}^{\prime} m_{2}^{\prime}\right\rangle\left\langle m_{1}^{\prime} m_{2}^{\prime} \mid j m\right\rangle \\
& =\hbar \sum_{m_{1}^{\prime} m_{2}^{\prime}}\left(m_{1}^{\prime}+m_{2}^{\prime}\right)\left|m_{1}^{\prime} m_{2}^{\prime}\right\rangle\left\langle m_{1}^{\prime} m_{2}^{\prime} \mid j m\right\rangle \tag{H.5}
\end{align*}
$$

We can also act directly on $|j m\rangle$ and then insert a complete set of states

$$
\begin{align*}
J_{z}|j m\rangle & =\sum_{m_{1}^{\prime} m_{2}^{\prime}}\left|m_{1}^{\prime} m_{2}^{\prime}\right\rangle\left\langle m_{1}^{\prime} m_{2}^{\prime}\right| J_{z}|j m\rangle \\
& =\hbar m \sum_{m_{1}^{\prime} m_{2}^{\prime}}\left|m_{1}^{\prime} m_{2}^{\prime}\right\rangle\left\langle m_{1}^{\prime} m_{2}^{\prime} \mid j m\right\rangle . \tag{H.6}
\end{align*}
$$

This provides us with the identity

$$
\begin{equation*}
m \sum_{m_{1}^{\prime} m_{2}^{\prime}}\left|m_{1}^{\prime} m_{2}^{\prime}\right\rangle\left\langle m_{1}^{\prime} m_{2}^{\prime} \mid j m\right\rangle=\sum_{m_{1}^{\prime} m_{2}^{\prime}}\left(m_{1}^{\prime}+m_{2}^{\prime}\right)\left|m_{1}^{\prime} m_{2}^{\prime}\right\rangle\left\langle m_{1}^{\prime} m_{2}^{\prime} \mid j m\right\rangle \tag{H.7}
\end{equation*}
$$

This equality must be valid for any $|j m\rangle$, and since all the kets $\left|m_{1}^{\prime} m_{2}^{\prime}\right\rangle$ are linearly independent, we must have for any $m_{1}^{\prime}, m_{2}^{\prime}$

$$
\begin{equation*}
\left(m-m_{1}^{\prime}-m_{2}^{\prime}\right)\left\langle m_{1}^{\prime} m_{2}^{\prime} \mid j m\right\rangle\left|m_{1}^{\prime} m_{2}^{\prime}\right\rangle=0 . \tag{H.8}
\end{equation*}
$$

We have two ways to get this zero. One of them is a $m=m_{1}^{\prime}+m_{2}^{\prime}$ condition, and the other is for the CG coeff $\left\langle m_{1}^{\prime} m_{2}^{\prime} \mid j m\right\rangle$ to be zero whenever $m \neq m_{1}^{\prime}+m_{2}^{\prime}$. It is not a difficult argument, but one that was not clear from a read of the text (at least to me).

ONE MORE ADIABATIC PERTURBATION I DERIVATION.

## I. 1 motivation.

I liked one of the adiabatic perturbation derivations that I did to review the material, and am recording it for reference.

## I. 2 BUILD UP.

In time dependent perturbation we started after noting that our ket in the interaction picture, for a Hamiltonian $H=H_{0}+H^{\prime}(t)$, took the form

$$
\begin{equation*}
\left|\alpha_{S}(t)\right\rangle=e^{-i H_{0} t / \hbar}\left|\alpha_{I}(t)\right\rangle=e^{-i H_{0} t / \hbar} U_{I}(t)\left|\alpha_{I}(0)\right\rangle . \tag{I.1}
\end{equation*}
$$

Here we have basically assumed that the time evolution can be factored into a portion dependent on only the static portion of the Hamiltonian, with some other operator $U_{I}(t)$, providing the remainder of the time evolution. From eq. (I.1) that operator $U_{I}(t)$ is found to behave according to

$$
\begin{equation*}
i \hbar \frac{d U_{I}}{d t}=e^{i H_{0} t / \hbar} H^{\prime}(t) e^{-i H_{0} t / \hbar} U_{I} \tag{I.2}
\end{equation*}
$$

but for our purposes we just assumed it existed, and used this for motivation. With the assumption that the interaction picture kets can be written in terms of the basis kets for the system at $t=0$ we write our Schrödinger ket as

$$
\begin{equation*}
|\psi\rangle=\sum_{k} e^{-i H_{0} t / \hbar} a_{k}(t)|k\rangle=\sum_{k} e^{-i \omega_{k} t / \hbar} a_{k}(t)|k\rangle \tag{I.3}
\end{equation*}
$$

where $|k\rangle$ are the energy eigenkets for the initial time equation problem

$$
\begin{equation*}
H_{0}|k\rangle=E_{k}^{0}|k\rangle \tag{I.4}
\end{equation*}
$$

## I. 3 ADIABATIC CASE.

For the adiabatic problem, we assume the system is changing very slowly, as described by the instantaneous energy eigenkets

$$
\begin{equation*}
H(t)|k(t)\rangle=E_{k}(t)|k(t)\rangle . \tag{I.5}
\end{equation*}
$$

Can we assume a similar representation to eq. (I.3) above, but allow $|k\rangle$ to vary in time? This does not quite work since $|k(t)\rangle$ are no longer eigenkets of $H_{0}$

$$
\begin{equation*}
|\psi\rangle=\sum_{k} e^{-i H_{0} t / \hbar} a_{k}(t)|k(t)\rangle \neq \sum_{k} e^{-i \omega_{k} t} a_{k}(t)|k(t)\rangle . \tag{I.6}
\end{equation*}
$$

Operating with $e^{i H_{0} t / \hbar}$ does not give the proper time evolution of $|k(t)\rangle$, and we will in general have a more complex functional dependence in our evolution operator for each $|k(t)\rangle$. Instead of an $\omega_{k} t$ dependence in this time evolution operator let us assume we have some function $\alpha_{k}(t)$ to be determined, and can write our ket as

$$
\begin{equation*}
|\psi\rangle=\sum_{k} e^{-i \alpha_{k}(t)} a_{k}(t)|k(t)\rangle \tag{I.7}
\end{equation*}
$$

Operating on this with our energy operator equation we have

$$
\begin{align*}
0 & =\left(H-i \hbar \frac{d}{d t}\right)|\psi\rangle \\
& =\left(H-i \hbar \frac{d}{d t}\right) \sum_{k} e^{-i \alpha_{k}} a_{k}|k\rangle  \tag{I.8}\\
& =\sum_{k} e^{-i \alpha_{k}(t)}\left(\left(E_{k} a_{k}-i \hbar\left(-i \alpha_{k}^{\prime} a_{k}+a_{k}^{\prime}\right)\right)|k\rangle-i \hbar a_{k}\left|k^{\prime}\right\rangle\right) .
\end{align*}
$$

Here I have written $\left|k^{\prime}\right\rangle=d|k\rangle / d t$. In our original time dependent perturbation the $-i \alpha_{k}^{\prime}$ term was $-i \omega_{k}$, so this killed off the $E_{k}$. If we assume this still kills off the $E_{k}$, we must have

$$
\begin{equation*}
\alpha_{k}=\frac{1}{\hbar} \int_{0}^{t} E_{k}\left(t^{\prime}\right) d t^{\prime}, \tag{I.9}
\end{equation*}
$$

and are left with

$$
\begin{equation*}
0=\sum_{k} e^{-i \alpha_{k}(t)}\left(a_{k}^{\prime}|k\rangle+a_{k}\left|k^{\prime}\right\rangle\right) . \tag{I.10}
\end{equation*}
$$

Bra'ing with $\langle m|$ we have

$$
\begin{equation*}
0=e^{-i \alpha_{m}(t)} a_{m}^{\prime}+e^{-i \alpha_{m}(t)} a_{m}\left\langle m \mid m^{\prime}\right\rangle+\sum_{k \neq m} e^{-i \alpha_{k}(t)} a_{k}\left\langle m \mid k^{\prime}\right\rangle \tag{I.11}
\end{equation*}
$$

or

$$
\begin{equation*}
a_{m}^{\prime}+a_{m}\left\langle m \mid m^{\prime}\right\rangle=-\sum_{k \neq m} e^{-i \alpha_{k}(t)} e^{i \alpha_{m}(t)} a_{k}\left\langle m \mid k^{\prime}\right\rangle \tag{I.12}
\end{equation*}
$$

The LHS is a perfect differential if we introduce an integration factor $e^{\int_{0}^{t}\left\langle m \mid m^{\prime}\right\rangle}$, so we can write

$$
\begin{equation*}
e^{-\int_{0}^{t}\left\langle m \mid m^{\prime}\right\rangle}\left(a_{m} e^{\int_{0}^{t}\left\langle m \mid m^{\prime}\right\rangle}\right)^{\prime}=-\sum_{k \neq m} e^{-i \alpha_{k}(t)} e^{i \alpha_{m}(t)} a_{k}\left\langle m \mid k^{\prime}\right\rangle \tag{I.13}
\end{equation*}
$$

This suggests that we want to form a new function

$$
\begin{equation*}
b_{m}=a_{m} e^{\int_{0}^{t}\left\langle m \mid m^{\prime}\right\rangle}, \tag{I.14}
\end{equation*}
$$

or

$$
\begin{equation*}
a_{m}=b_{m} e^{-\int_{0}^{t}\left\langle m \mid m^{\prime}\right\rangle} \tag{I.15}
\end{equation*}
$$

Plugging this into our assumed representation we have a more concrete form

$$
\begin{equation*}
|\psi\rangle=\sum_{k} e^{-\int_{0}^{t} d t^{\prime}\left(i \omega_{k}+\left\langle k \mid k^{\prime}\right\rangle\right)} b_{k}(t)|k(t)\rangle \tag{I.16}
\end{equation*}
$$

Writing

$$
\begin{equation*}
\Gamma_{k}=i\left\langle k \mid k^{\prime}\right\rangle \tag{I.17}
\end{equation*}
$$

this becomes

$$
\begin{equation*}
|\psi\rangle=\sum_{k} e^{-i \int_{0}^{t} d t^{\prime}\left(\omega_{k}-\Gamma_{k}\right)} b_{k}(t)|k(t)\rangle \tag{I.18}
\end{equation*}
$$

A final pass. Now that we have what appears to be a good representation for any given state if we wish to examine the time evolution, let us start over, reapplying our instantaneous energy operator equality

$$
\begin{align*}
0 & =\left(H-i \hbar \frac{d}{d t}\right)|\psi\rangle \\
& =\left(H-i \hbar \frac{d}{d t}\right) \sum_{k} e^{-i \int_{0}^{t} d t^{\prime}\left(\omega_{k}-\Gamma_{k}\right)} b_{k}|k\rangle  \tag{I.19}\\
& =-i \hbar \sum_{k} e^{-i \int_{0}^{t} d t^{\prime}\left(\omega_{k}-\Gamma_{k}\right)}\left(i \Gamma_{k} b_{k}|k\rangle+b_{k}^{\prime}|k\rangle+b_{k}\left|k^{\prime}\right\rangle\right)
\end{align*}
$$

Bra'ing with $\langle m|$ we find

$$
\begin{align*}
0 & =e^{-i \int_{0}^{t} d t^{\prime}\left(\omega_{m}-\Gamma_{m}\right)} i \Gamma_{m} b_{m}+e^{-i \int_{0}^{t} d t^{\prime}\left(\omega_{m}-\Gamma_{m}\right)} b_{m}^{\prime} \\
& +e^{-i \int_{0}^{t} d t^{\prime}\left(\omega_{m}-\Gamma_{m}\right)} b_{m}\left\langle m \mid m^{\prime}\right\rangle+\sum_{k \neq m} e^{-i \int_{0}^{t} d t^{\prime}\left(\omega_{k}-\Gamma_{k}\right)} b_{k}\left\langle m \mid k^{\prime}\right\rangle . \tag{I.20}
\end{align*}
$$

Since $i \Gamma_{m}=\left\langle m \mid m^{\prime}\right\rangle$ the first and third terms cancel leaving us just

$$
\begin{equation*}
b_{m}^{\prime}=-\sum_{k \neq m} e^{-i \int_{0}^{t} d t^{\prime}\left(\omega_{k m}-\Gamma_{k m}\right)} b_{k}\left\langle m \mid k^{\prime}\right\rangle \tag{I.21}
\end{equation*}
$$

where $\omega_{k m}=\omega_{k}-\omega_{m}$ and $\Gamma_{k m}=\Gamma_{k}-\Gamma_{m}$.
I. 4 SUMMARY.

We assumed that a ket for the system has a representation in the form

$$
\begin{equation*}
|\psi\rangle=\sum_{k} e^{-i \alpha_{k}(t)} a_{k}(t)|k(t)\rangle \tag{I.22}
\end{equation*}
$$

where $a_{k}(t)$ and $\alpha_{k}(t)$ are given or to be determined. Application of our energy operator identity provides us with an alternate representation that simplifies the results

$$
\begin{equation*}
|\psi\rangle=\sum_{k} e^{-i \int_{0}^{t} d t^{\prime}\left(\omega_{k}-\Gamma_{k}\right)} b_{k}(t)|k(t)\rangle \tag{I.23}
\end{equation*}
$$

With

$$
\begin{align*}
\left|m^{\prime}\right\rangle & =\frac{d}{d t}|m\rangle \\
\Gamma_{k} & =i\left\langle m \mid m^{\prime}\right\rangle  \tag{I.24}\\
\omega_{k m} & =\omega_{k}-\omega_{m} \\
\Gamma_{k m} & =\Gamma_{k}-\Gamma_{m},
\end{align*}
$$

we find that our dynamics of the coefficients are related by

$$
\begin{equation*}
b_{m}^{\prime}=-\sum_{k \neq m} e^{-i \int_{0}^{t} d t^{\prime}\left(\omega_{k m}-\Gamma_{k m}\right)} b_{k}\left\langle m \mid k^{\prime}\right\rangle \tag{I.25}
\end{equation*}
$$

## TIME DEPENDENT PERTURBATION REVISITED.

With

$$
\begin{equation*}
|\psi t\rangle=\sum_{k} c_{k}(t) e^{-i \omega_{k} t}|k\rangle \tag{J.1}
\end{equation*}
$$

apply the energy eigenvalue operator identity

$$
\begin{align*}
0 & =\left(H_{0}+H^{\prime}-i \hbar \frac{d}{d t}\right)|\psi t\rangle \\
& =\left(H_{0}+H^{\prime}-i \hbar \frac{d}{d t}\right) \sum_{k} c_{k} e^{-i \omega_{k} t}|k\rangle  \tag{J.2}\\
& =\sum_{k} e^{-i \omega_{k} t}\left(c_{k} E_{k}+H^{\prime} c_{k}-i \hbar\left(-i \omega_{k}\right) c_{k}-i \hbar c_{k}^{\prime}\right)|k\rangle
\end{align*}
$$

Bra with $\langle m|$

$$
\begin{equation*}
\sum_{k} e^{-i \omega_{k} t} H_{m k}^{\prime} c_{k}=i \hbar e^{-i \omega_{m} t} c_{m}^{\prime} \tag{J.3}
\end{equation*}
$$

or

$$
\begin{equation*}
c_{m}^{\prime}=\frac{1}{i \hbar} \sum_{k} e^{-i \omega_{k m} t} H_{m k}^{\prime} c_{k} \tag{J.4}
\end{equation*}
$$

Now we can make the assumptions about the initial state and away we go.

## SECOND FORM OF ADIABATIC APPROXIMATION.

Motivation. In class we were shown an adiabatic approximation where we started with (or worked our way towards) a representation of the form

$$
\begin{equation*}
|\psi\rangle=\sum_{k} c_{k}(t) e^{-i} \int_{0}^{t}\left(\omega_{k}\left(t^{\prime}\right)-\Gamma_{k}\left(t^{\prime}\right)\right) d t^{\prime}\left|\psi_{k}(t)\right\rangle, \tag{K.1}
\end{equation*}
$$

where $\left|\psi_{k}(t)\right\rangle$ were normalized energy eigenkets for the (slowly) evolving Hamiltonian

$$
\begin{equation*}
H(t)\left|\psi_{k}(t)\right\rangle=E_{k}(t)\left|\psi_{k}(t)\right\rangle . \tag{K.2}
\end{equation*}
$$

In the problem sets we were shown a different adiabatic approximation, where are starting point is

$$
|\psi(t)\rangle=\sum_{k} c_{k}(t)\left|\psi_{k}(t)\right\rangle .
$$

For completeness, here is a walk through of the general amplitude derivation that is been used.

Guts. We operate with our energy identity once again

$$
\begin{align*}
0 & =\left(H-i \hbar \frac{d}{d t}\right) \sum_{k} c_{k}|k\rangle  \tag{K.4}\\
& =\sum_{k} c_{k} E_{k}|k\rangle-i \hbar c_{k}^{\prime}|k\rangle-i \hbar c_{k}\left|k^{\prime}\right\rangle,
\end{align*}
$$

where

$$
\begin{equation*}
\left|k^{\prime}\right\rangle=\frac{d}{d t}|k\rangle . \tag{K.5}
\end{equation*}
$$

Bra'ing with $\langle m|$, and split the sum into $k=m$ and $k \neq m$ parts

$$
\begin{equation*}
0=c_{m} E_{m}-i \hbar c_{m}^{\prime}-i \hbar c_{m}\left\langle m \mid m^{\prime}\right\rangle-i \hbar \sum_{k \neq m} c_{k}\left\langle m \mid k^{\prime}\right\rangle \tag{K.6}
\end{equation*}
$$

Again writing

$$
\begin{equation*}
\Gamma_{m}=i\left\langle m \mid m^{\prime}\right\rangle \tag{K.7}
\end{equation*}
$$

We have

$$
\begin{equation*}
c_{m}^{\prime}=\frac{1}{i \hbar} c_{m}\left(E_{m}-\hbar \Gamma_{m}\right)-\sum_{k \neq m} c_{k}\left\langle m \mid k^{\prime}\right\rangle \tag{K.8}
\end{equation*}
$$

In this form we can make an "Adiabatic" approximation, dropping the $k \neq m$ terms, and integrate

$$
\begin{equation*}
\int \frac{d c_{m}^{\prime}}{c_{m}}=\frac{1}{i \hbar} \int_{0}^{t}\left(E_{m}\left(t^{\prime}\right)-\hbar \Gamma_{m}\left(t^{\prime}\right)\right) d t^{\prime} \tag{K.9}
\end{equation*}
$$

or

$$
\begin{equation*}
c_{m}(t)=A \exp \left(\frac{1}{i \hbar} \int_{0}^{t}\left(E_{m}\left(t^{\prime}\right)-\hbar \Gamma_{m}\left(t^{\prime}\right)\right) d t^{\prime}\right) \tag{K.10}
\end{equation*}
$$

Evaluating at $t=0$, fixes the integration constant for

$$
\begin{equation*}
c_{m}(t)=c_{m}(0) \exp \left(\frac{1}{i \hbar} \int_{0}^{t}\left(E_{m}\left(t^{\prime}\right)-\hbar \Gamma_{m}\left(t^{\prime}\right)\right) d t^{\prime}\right) \tag{K.11}
\end{equation*}
$$

Observe that this is very close to the starting point of the adiabatic approximation we performed in class since we end up with

$$
\begin{equation*}
|\psi\rangle=\sum_{k} c_{k}(0) e^{-i \int_{0}^{t}\left(\omega_{k}\left(t^{\prime}\right)-\Gamma_{k}\left(t^{\prime}\right)\right) d t^{\prime}}|k(t)\rangle . \tag{K.12}
\end{equation*}
$$

So, to perform the more detailed approximation, that started with eq. (K.1), where we ended up with all the cross terms that had both $\omega_{k}$ and Berry phase $\Gamma_{k}$ dependence, we have only to generalize by replacing $c_{k}(0)$ with $c_{k}(t)$.

Motivation. In class this week, looking at an instance of the Helmholtz equation

$$
\begin{equation*}
\left(\boldsymbol{\nabla}^{2}+\mathbf{k}^{2}\right) \psi_{\mathbf{k}}(\mathbf{r})=s(\mathbf{r}) \tag{L.1}
\end{equation*}
$$

We were told that the Green's function

$$
\begin{equation*}
\left(\boldsymbol{\nabla}^{2}+\mathbf{k}^{2}\right) G^{0}\left(\mathbf{r}, \mathbf{r}^{\prime}\right)=\delta\left(\mathbf{r}-\mathbf{r}^{\prime}\right) \tag{L.2}
\end{equation*}
$$

that can be used to solve for a particular solution this differential equation via convolution

$$
\begin{equation*}
\psi_{\mathbf{k}}(\mathbf{r})=\int G^{0}\left(\mathbf{r}, \mathbf{r}^{\prime}\right) s\left(\mathbf{r}^{\prime}\right) d^{3} \mathbf{r}^{\prime} \tag{L.3}
\end{equation*}
$$

had the value

$$
\begin{equation*}
G^{0}\left(\mathbf{r}, \mathbf{r}^{\prime}\right)=-\frac{1}{4 \pi} \frac{e^{i k\left|\mathbf{r}-\mathbf{r}^{\prime}\right|}}{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|} \tag{L.4}
\end{equation*}
$$

Let us try to verify this. Application of the Helmholtz differential operator $\boldsymbol{\nabla}^{2}+\mathbf{k}^{2}$ on the presumed solution gives

$$
\begin{equation*}
\left(\boldsymbol{\nabla}^{2}+\mathbf{k}^{2}\right) \psi_{\mathbf{k}}(\mathbf{r})=-\frac{1}{4 \pi} \int\left(\boldsymbol{\nabla}^{2}+\mathbf{k}^{2}\right) \frac{e^{i k\left|\mathbf{r}-\mathbf{r}^{\prime}\right|}}{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|} s\left(\mathbf{r}^{\prime}\right) d^{3} \mathbf{r}^{\prime} \tag{L.5}
\end{equation*}
$$

When $\mathbf{r} \neq \mathbf{r}^{\prime}$. To proceed we will need to evaluate

$$
\begin{equation*}
\nabla^{2} \frac{e^{i k\left|\mathbf{r}-\mathbf{r}^{\prime}\right|}}{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|} \tag{L.6}
\end{equation*}
$$

Writing $\mu=\left|\mathbf{r}-\mathbf{r}^{\prime}\right|$ we start with the computation of

$$
\begin{align*}
\frac{\partial}{\partial x} \frac{e^{i k \mu}}{\mu} & =\frac{\partial \mu}{\partial x}\left(\frac{i k}{\mu}-\frac{1}{\mu^{2}}\right) e^{i k \mu}  \tag{L.7}\\
& =\frac{\partial \mu}{\partial x}\left(i k-\frac{1}{\mu}\right) \frac{e^{i k \mu}}{\mu}
\end{align*}
$$

We see that we will have

$$
\begin{equation*}
\nabla \frac{e^{i k \mu}}{\mu}=\left(i k-\frac{1}{\mu}\right) \frac{e^{i k \mu}}{\mu} \nabla \mu \tag{L.8}
\end{equation*}
$$

Taking second derivatives with respect to $x$ we find

$$
\begin{align*}
\frac{\partial^{2}}{\partial x^{2}} \frac{e^{i k \mu}}{\mu} & =\frac{\partial^{2} \mu}{\partial x^{2}}\left(i k-\frac{1}{\mu}\right) \frac{e^{i k \mu}}{\mu}+\frac{\partial \mu}{\partial x} \frac{\partial \mu}{\partial x} \frac{1}{\mu^{2}} \frac{e^{i k \mu}}{\mu}+\left(\frac{\partial \mu}{\partial x}\right)^{2}\left(i k-\frac{1}{\mu}\right)^{2} \frac{e^{i k \mu}}{\mu} \\
& =\frac{\partial^{2} \mu}{\partial x^{2}}\left(i k-\frac{1}{\mu}\right) \frac{e^{i k \mu}}{\mu}+\left(\frac{\partial \mu}{\partial x}\right)^{2}\left(-k^{2}-\frac{2 i k}{\mu}+\frac{2}{\mu^{2}}\right) \frac{e^{i k \mu}}{\mu} . \tag{L.9}
\end{align*}
$$

Our Laplacian is then

$$
\begin{equation*}
\nabla^{2} \frac{e^{i k \mu}}{\mu}=\left(i k-\frac{1}{\mu}\right) \frac{e^{i k \mu}}{\mu} \nabla^{2} \mu+\left(-k^{2}-\frac{2 i k}{\mu}+\frac{2}{\mu^{2}}\right) \frac{e^{i k \mu}}{\mu}(\nabla \mu)^{2} . \tag{L.10}
\end{equation*}
$$

Now lets calculate the derivatives of $\mu$. Working on $x$ again, we have

$$
\begin{align*}
\frac{\partial}{\partial x} \mu & =\frac{\partial}{\partial x} \sqrt{\left(x-x^{\prime}\right)^{2}+\left(y-y^{\prime}\right)^{2}+\left(z-z^{\prime}\right)^{2}} \\
& =\frac{1}{2} 2\left(x-x^{\prime}\right) \frac{1}{\sqrt{\left(x-x^{\prime}\right)^{2}+\left(y-y^{\prime}\right)^{2}+\left(z-z^{\prime}\right)^{2}}}  \tag{L.11}\\
& =\frac{x-x^{\prime}}{\mu} .
\end{align*}
$$

So we have

$$
\begin{align*}
\boldsymbol{\nabla} \mu & =\frac{\mathbf{r}-\mathbf{r}^{\prime}}{\mu}  \tag{L.12}\\
(\nabla \mu)^{2} & =1 .
\end{align*}
$$

Taking second derivatives with respect to $x$ we find

$$
\begin{align*}
\frac{\partial^{2}}{\partial x^{2}} \mu & =\frac{\partial}{\partial x} \frac{x-x^{\prime}}{\mu} \\
& =\frac{1}{\mu}-\left(x-x^{\prime}\right) \frac{\partial \mu}{\partial x} \frac{1}{\mu^{2}}  \tag{L.13}\\
& =\frac{1}{\mu}-\left(x-x^{\prime}\right) \frac{x-x^{\prime}}{\mu} \frac{1}{\mu^{2}} \\
& =\frac{1}{\mu}-\left(x-x^{\prime}\right)^{2} \frac{1}{\mu^{3}} .
\end{align*}
$$

So we find

$$
\begin{equation*}
\nabla^{2} \mu=\frac{3}{\mu}-\frac{1}{\mu}, \tag{L.14}
\end{equation*}
$$

or

$$
\begin{equation*}
\nabla^{2} \mu=\frac{2}{\mu} \tag{L.15}
\end{equation*}
$$

Inserting this and $(\nabla \mu)^{2}$ into eq. (L.10) we find

$$
\begin{equation*}
\nabla^{2} \frac{e^{i k \mu}}{\mu}=\left(i k-\frac{1}{\mu}\right) \frac{e^{i k \mu}}{\mu} \frac{2}{\mu}+\left(-k^{2}-\frac{2 i k}{\mu}+\frac{2}{\mu^{2}}\right) \frac{e^{i k \mu}}{\mu}=-k^{2} \frac{e^{i k \mu}}{\mu} . \tag{L.16}
\end{equation*}
$$

This shows us that provided $\mathbf{r} \neq \mathbf{r}^{\prime}$ we have

$$
\begin{equation*}
\left(\boldsymbol{\nabla}^{2}+\mathbf{k}^{2}\right) G^{0}\left(\mathbf{r}, \mathbf{r}^{\prime}\right)=0 \tag{L.17}
\end{equation*}
$$

In the neighborhood of $\left|\mathbf{r}-\mathbf{r}^{\prime}\right|<\epsilon$. Having shown that we end up with zero everywhere that $\mathbf{r} \neq \mathbf{r}^{\prime}$ we are left to consider a neighborhood of the volume surrounding the point $\mathbf{r}$ in our integral. Following the Coulomb treatment in $\S 2.2$ of [14] we use a spherical volume element centered around $\mathbf{r}$ of radius $\epsilon$, and then convert a divergence to a surface area to evaluate the integral away from the problematic point

$$
\begin{align*}
- & \frac{1}{4 \pi} \int_{\text {all space }}\left(\nabla^{2}+\mathbf{k}^{2}\right) \frac{e^{i k\left|\mathbf{r}-\mathbf{r}^{\prime}\right|}}{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|} s\left(\mathbf{r}^{\prime}\right) d^{3} \mathbf{r}^{\prime}  \tag{L.18}\\
& =-\frac{1}{4 \pi} \int_{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|<\epsilon}\left(\nabla^{2}+\mathbf{k}^{2}\right) \frac{e^{i k\left|\mathbf{r}-\mathbf{r}^{\prime}\right|}}{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|} s\left(\mathbf{r}^{\prime}\right) d^{3} \mathbf{r}^{\prime} .
\end{align*}
$$

We make the change of variables $\mathbf{r}^{\prime}=\mathbf{r}+\mathbf{a}$. We add an explicit $\mathbf{r}$ suffix to our Laplacian at the same time to remind us that it is taking derivatives with respect to the coordinates of $\mathbf{r}=(x, y, z)$, and not the coordinates of our integration variable $\mathbf{a}=\left(a_{x}, a_{y}, a_{z}\right)$. Assuming sufficient continuity and "well behavedness" of $s\left(\mathbf{r}^{\prime}\right)$ we will be able to pull it out of the integral, giving

$$
\begin{align*}
-\frac{1}{4 \pi} & \int_{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|<\epsilon}\left(\nabla_{\mathbf{r}}^{2}+\mathbf{k}^{2}\right) \frac{e^{i k\left|\mathbf{r}-\mathbf{r}^{\prime}\right|}}{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|} s\left(\mathbf{r}^{\prime}\right) d^{3} \mathbf{r}^{\prime} \\
& =-\frac{1}{4 \pi} \int_{|\mathbf{a}|<\epsilon}\left(\nabla_{\mathbf{r}}^{2}+\mathbf{k}^{2}\right) \frac{e^{i k|\mathbf{a}|}}{|\mathbf{a}|} s(\mathbf{r}+\mathbf{a}) d^{3} \mathbf{a}  \tag{L.19}\\
& =-\frac{s(\mathbf{r})}{4 \pi} \int_{|\mathbf{a}|<\epsilon}\left(\nabla_{\mathbf{r}}^{2}+\mathbf{k}^{2}\right) \frac{e^{i k|\mathbf{a}|}}{|\mathbf{a}|} d^{3} \mathbf{a}
\end{align*}
$$

Recalling the dependencies on the derivatives of $\left|\mathbf{r}-\mathbf{r}^{\prime}\right|$ in our previous gradient evaluations, we note that we have

$$
\begin{align*}
\boldsymbol{\nabla}_{\mathbf{r}}\left|\mathbf{r}-\mathbf{r}^{\prime}\right| & =-\boldsymbol{\nabla}_{\mathbf{a}}|\mathbf{a}| \\
\left(\boldsymbol{\nabla}_{\mathbf{r}}\left|\mathbf{r}-\mathbf{r}^{\prime}\right|\right)^{2} & =\left(\boldsymbol{\nabla}_{\mathbf{a}}|\mathbf{a}|\right)^{2}  \tag{L.20}\\
\boldsymbol{\nabla}_{\mathbf{r}}^{2}\left|\mathbf{r}-\mathbf{r}^{\prime}\right| & =\boldsymbol{\nabla}_{\mathbf{a}}^{2}|\mathbf{a}|,
\end{align*}
$$

so with $\mathbf{a}=\mathbf{r}-\mathbf{r}^{\prime}$, we can rewrite our Laplacian as

$$
\begin{equation*}
\nabla_{\mathbf{r}}^{2} \frac{e^{i k\left|\mathbf{r}-\mathbf{r}^{\prime}\right|}}{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|}=\nabla_{\mathbf{a}}^{2} \frac{e^{i k|\mathbf{a}|}}{|\mathbf{a}|}=\boldsymbol{\nabla}_{\mathbf{a}} \cdot\left(\boldsymbol{\nabla}_{\mathbf{a}} \frac{e^{i k|\mathbf{a}|}}{|\mathbf{a}|}\right) \tag{L.21}
\end{equation*}
$$

This gives us

$$
\begin{align*}
-\frac{s(\mathbf{r})}{4 \pi} & \int_{|\mathbf{a}|<\epsilon}\left(\nabla_{\mathbf{a}}^{2}+\mathbf{k}^{2}\right) \frac{e^{i k|\mathbf{a}|}}{|\mathbf{a}|} d^{3} \mathbf{a} \\
& =-\frac{s(\mathbf{r})}{4 \pi} \int_{d V} \boldsymbol{\nabla}_{\mathbf{a}} \cdot\left(\boldsymbol{\nabla}_{\mathbf{a}} \frac{e^{i k|\mathbf{a}|}}{|\mathbf{a}|}\right) d^{3} \mathbf{a}-\frac{s(\mathbf{r})}{4 \pi} \int_{d V} \mathbf{k}^{2} \frac{e^{i k|\mathbf{a}|}}{|\mathbf{a}|} d^{3} \mathbf{a}  \tag{L.22}\\
& =-\frac{s(\mathbf{r})}{4 \pi} \int_{d A}\left(\boldsymbol{\nabla}_{\mathbf{a}} \frac{e^{i k|\mathbf{a}|}}{|\mathbf{a}|}\right) \cdot \hat{\mathbf{a}} d^{2} \mathbf{a}-\frac{s(\mathbf{r})}{4 \pi} \int_{d V} \mathbf{k}^{2} \frac{e^{i k|\mathbf{a}|}}{|\mathbf{a}|} d^{3} \mathbf{a}
\end{align*}
$$

To complete these evaluations, we can now employ a spherical coordinate change of variables. Let us do the $\mathbf{k}^{2}$ volume integral first. We have

$$
\begin{align*}
\int_{d V} \mathbf{k}^{2} \frac{e^{i k|\mathbf{a}|}}{|\mathbf{a}|} d^{3} \mathbf{a} & =\int_{a=0}^{\epsilon} \int_{\theta=0}^{\pi} \int_{\phi=0}^{2 \pi} \mathbf{k}^{2} \frac{e^{i k a}}{a} a^{2} d a \sin \theta d \theta d \phi \\
& =4 \pi k^{2} \int_{a=0}^{\epsilon} a e^{i k a} d a \\
& =4 \pi \int_{u=0}^{k \epsilon} u e^{i u} d u  \tag{L.23}\\
& =\left.4 \pi(-i u+1) e^{i u}\right|_{0} ^{k \epsilon} \\
& =4 \pi\left((-i k \epsilon+1) e^{i k \epsilon}-1\right)
\end{align*}
$$

To evaluate the surface integral we note that we will require only the radial portion of the gradient, so have

$$
\begin{align*}
\left(\boldsymbol{\nabla}_{\mathbf{a}} \frac{e^{i k|\mathbf{a}|}}{|\mathbf{a}|}\right) \cdot \hat{\mathbf{a}} & =\left(\hat{\mathbf{a}} \frac{\partial}{\partial a} \frac{e^{i k a}}{a}\right) \cdot \hat{\mathbf{a}} \\
& =\frac{\partial}{\partial a} \frac{e^{i k a}}{a}  \tag{L.24}\\
& =\left(i k \frac{1}{a}-\frac{1}{a^{2}}\right) e^{i k a} \\
& =(i k a-1) \frac{e^{i k a}}{a^{2}}
\end{align*}
$$

Our area element is $a^{2} \sin \theta d \theta d \phi$, so we are left with

$$
\begin{align*}
\int_{d A}\left(\nabla_{\mathbf{a}} \frac{e^{i k|\mathbf{a}|}}{|\mathbf{a}|}\right) \cdot \hat{\mathbf{a}} d^{2} \mathbf{a} & =\left.\int_{\theta=0}^{\pi} \int_{\phi=0}^{2 \pi}(i k a-1) \frac{e^{i k a}}{a^{2}} a^{2} \sin \theta d \theta d \phi\right|_{a=\epsilon}  \tag{L.25}\\
& =4 \pi(i k \epsilon-1) e^{i k \epsilon}
\end{align*}
$$

Putting everything back together we have

$$
\begin{align*}
-\frac{1}{4 \pi} & \int_{\text {all space }}\left(\nabla^{2}+\mathbf{k}^{2}\right) \frac{e^{i k\left|\mathbf{r}-\mathbf{r}^{\prime}\right|}}{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|} s\left(\mathbf{r}^{\prime}\right) d^{3} \mathbf{r}^{\prime} \\
& =-s(\mathbf{r})\left((-i k \epsilon+1) e^{i k \epsilon}-1+(i k \epsilon-1) e^{i k \epsilon}\right)  \tag{L.26}\\
& =-s(\mathbf{r})\left((-i k \epsilon+1+i k \epsilon-1) e^{i k \epsilon}-1\right) .
\end{align*}
$$

This is just

$$
\begin{equation*}
-\frac{1}{4 \pi} \int_{\text {all space }}\left(\nabla^{2}+\mathbf{k}^{2}\right) \frac{e^{i k\left|\mathbf{r}-\mathbf{r}^{\prime}\right|}}{\left|\mathbf{r}-\mathbf{r}^{\prime}\right|} s\left(\mathbf{r}^{\prime}\right) d^{3} \mathbf{r}^{\prime}=s(\mathbf{r}) \tag{L.27}
\end{equation*}
$$

This completes the desired verification of the Green's function for the Helmholtz operator. Observe the perfect cancellation here, so the limit of $\epsilon \rightarrow 0$ can be independent of how large $k$ is made. You have to complete the integrals for both the Laplacian and the $\mathbf{k}^{2}$ portions of the integrals and add them, before taking any limits, or else you will get into trouble (as I did in my first attempt).

These Mathematica notebooks, some just trivial ones used to generate figures, others more elaborate, and perhaps some even polished, can be found in
https://github.com/peeterjoot/mathematica/tree/master/phy456/.
The free Wolfram CDF player, is capable of read-only viewing these notebooks to some extent.

Files saved explicitly as CDF have interactive content that can be explored with the CDF player.

- Sep 15, 2011 desai_S_24_2_1_verify.nb

Some integrals related to QM hydrogen atom energy expectation values.

- Sep 24, 2011 ps2_p2_verify_normalization.nb

Some trig integrals that I didn't feel like doing manually.

- Sep 24, 2011 exponential_integrals.nb

More Gaussian integrals and some that Mathematica didn't know how to do.

- Sep 28, 2011 problem_set_3_integrals.nb

Some Gaussian integrals.

- Oct 2, 2011 24.4.3_attempt_with_mathematica.nb

Some variational method calculations for QM energy estimation.

- Oct 5, 2011 gaussian_fitting_for_abs_function.nb

Hankle function fitting for $\mathrm{e}^{\wedge}-\mathrm{b}|\mathrm{x}|$ and related plots.

- Oct 6, 2011 stackoverflow_expNth_unknown.nb

Stripped down example notebook for stackoverflow question about Derivative[2] not behaving well.

- Oct 6, 2011 stackoverflow_question_about_listable.nb Stripped down example notebook for stackoverflow question about Listable attribute defaults.
- Oct 8, 2011 qmTwoL8figures.nb

Plot of Gaussian weighted cosine, its Fourier transform, and figure for perturbation of Harmonic oscillator system.

- Oct 9, 2011 qmTwoL9figures.nb

Sinusoid plot turned on at $\mathrm{t} \_0$ and ongoing from there.

- Oct 10, 2011 problem_set_4,_problem_2.nb

Some trig integrals that Mathematica didn't evaluate correctly. Don't trust a tool without thinking whether the results are good!

- Oct 15, 2011 desai_24_4_4.nb

Another worked variational method problem.

- Oct 15, 2011 desai_24_4_5.nb

Another worked variational method problem.

- Oct 15, 2011 desai_24_4_6.nb

Another worked variational method problem. Looks like I've learned about the /. operator for evaluating variables with values.

- Oct 15, 2011 qmTwoL10figures.nb

Some sinc function plots. Learned how to use Manipulate to make sliders.

- Oct 16, 2011 desai_attempt_to_verify_section_16.3.nb Some energy expectation value calculations.
- Oct 17, 2011 qmTwoL11figures.nb

Some vector addition and function translation figures.

- Oct 18, 2011 problem_set_5_integrals.nb

Some integrals of first order linear polynomials.

- Oct 19, 2011 qmTwoL12_figures.nb

Some step and rect function plots.

- Oct 31, 2011 plot_question.nb

Another stackoverflow mathematica question. Why no output in my plot. Learned about Mathematica local and global variables as a result.

- Oct 31, 2011 ps7_rotation_matrix_orthonormal.nb

A sanity check on a rotation matrix calculated as part of a problem set.

- Dec 17, 2011 qmTwoExamReflection.cdf

Exam problem 2a. Calculate the matrix of a Perturbation Hamiltonian $-\mu_{d} \cdot \mathbf{E}$ with respect to the $n=2$ hydrogen atom wave functions.

- March 28, 2013 24.4.3.newAttempt.nb

A new attempt at Desai 24.4.3 from scratch. This one has an error, as did the original. The original is now fixed.

- Apr 19, 2020 hydrogenGroundStateFig2.nb

Plot of the magnitude of the hydrogen ground state wavefunction.

- Apr 19, 2020 hydrogenGroundStateVariational.nb

An attempt to confirm the variational calculations for the Hydrogen ground state energy. I cannot reproduce what was shown in class.

- Apr 20, 2020 kineticEnergyGreaterThanPotentialFig9.nb

Plot of Kinetic Energy greater than various potentials (scattering).

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